Asymptotics and numerical efficiency of the Allen-Cahn model for phase interfaces with low energy in solids

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Abstract

We study how the propagation speed of interfaces in the Allen-Cahn phase field model for phase transformations in solids consisting of the elasticity equations and the Allen-Cahn equation depends on two parameters of the model. The two parameters control the interface energy and the interface width but change also the interface speed. To this end we derive an asymptotic expansion of second order for the interface speed, called the kinetic relation, and prove that it is uniformly valid in both parameters. As a consequence we show that the model error is proportional to the interface width divided by the interface energy. We conclude that simulations of interfaces with low interface energy based on this model require a very small interface width, implying a large numerical effort. Effective simulations thus need adaptive mesh refinement or other advanced techniques.

This version of the paper contains the proofs of Theorem 4.5 and Lemma 5.8, which are omitted in the version published in Continuum Mechanics and Thermodynamics.

Key words: Allen-Cahn phase field model for solids, asymptotic expansion, propagation speed of phase interfaces, kinetic relation, model error, numerical efficiency

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1 Introduction

In this paper we study how the propagation speed of interfaces in the Allen-Cahn phase field model for phase transformations in elastic solids depends on two parameters of the model. The model consists of the partial differential equations of linear elasticity coupled to the standard Allen-Cahn phase field equation. The two parameters, which we denote by $\mu$ and $\lambda$, control the interface energy and the interface width, but variation of these parameters also changes the interface speed, or more precisely, the form of the relation, which determines the interface speed as a function of the stress field and the curvature of the interface. In sharp interface models this relation is called kinetic relation. We use this notion also for phase field models. Our goal is therefore to determine the kinetic relation for the Allen-Cahn model and the dependence of it on the two model parameters. To this end we must derive an asymptotic expansion for the propagation speed of the interface.
and prove an error estimate for this asymptotic expansion, which holds uniformly in both parameters. Our results have consequences for the efficiency of the Allen-Cahn model in numerical simulations of interfaces with small interface energy. These consequences are also discussed.

Let $\Omega \subseteq \mathbb{R}^3$ be a bounded open set with a sufficiently smooth boundary $\partial \Omega$. The points of $\Omega$ represent the material points of a solid elastic body. The unknown functions in the model are the displacement $u(t, x) \in \mathbb{R}$ of the material point $x$ at time $t$, the Cauchy stress tensor $T(t, x) \in S^3$, where $S^3$ denotes the set of all symmetric $3 \times 3$-matrices, and the order parameter $S(t, x) \in \mathbb{R}$. These unknowns must satisfy the model equations

$$-\text{div}_x T = b,$$  \quad (1.1)  

$$T = D(\varepsilon(\nabla_x u) - \varepsilon S),$$  \quad (1.2)  

$$\partial_t S = -\frac{c}{(\mu \lambda)^{1/2}} \left( \partial_S W(\varepsilon(\nabla_x u), S) + \frac{1}{\mu^{1/2}} \tilde{\psi}'(S) - \mu^{1/2} \Delta_S S \right),$$  \quad (1.3)  

in the domain $[0, \infty) \times \Omega$. The boundary and initial conditions are

$$u(t, x) = U(t, x), \quad (t, x) \in [0, \infty) \times \partial \Omega, \quad (1.4)$$  

$$\partial_{n_{\partial \Omega}} S(t, x) = 0, \quad (t, x) \in [0, \infty) \times \partial \Omega, \quad (1.5)$$  

$$S(0, x) = S(x), \quad x \in \Omega. \quad (1.6)$$  

Here $b(t, x) \in \mathbb{R}^3$, $U(t, x) \in \mathbb{R}^3$, $S(t, x) \in \mathbb{R}$ denote given data, the volume force, boundary displacement and initial data. $\partial_{n_{\partial \Omega}}$ denotes the derivative in direction of the unit normal vector $n_{\partial \Omega}$ to the boundary. The deformation gradient $\nabla_x u(t, x)$ is the $3 \times 3$-matrix of first order partial derivatives of $u$ with respect to the components $x_k$ of $x$, and the strain tensor

$$\varepsilon(\nabla_x u) = \frac{1}{2} \left( \nabla_x u + (\nabla_x u)^T \right)$$

is the symmetric part of the deformation gradient, where $(\nabla_x u)^T$ denotes the transpose matrix. The elasticity tensor $D : S^3 \rightarrow S^3$ is a linear symmetric, positive definite mapping, $\varepsilon \in S^3$ is a given constant matrix, the transformation strain, and $\mu > 0$ and $\lambda > 0$ are parameters. The elastic energy is given by

$$W(\varepsilon(\nabla_x u), S) = \frac{1}{2} \left( D(\varepsilon(\nabla_x u) - \varepsilon S) : (\varepsilon(\nabla_x u) - \varepsilon S) \right),$$  \quad (1.7)  

with the matrix scalar product $A : B = \sum_{i,j} a_{ij} b_{ij}$. Using (1.2), we obtain for the derivative

$$\partial_S W(\varepsilon, S) = -\varepsilon : D(\varepsilon(\nabla_x u) - \varepsilon S) = -\varepsilon : T.$$  \quad (1.8)  

$c > 0$ is a given constant and $\hat{\psi} : \mathbb{R} \rightarrow [0, \infty)$ is a double well potential satisfying

$$\hat{\psi}(0) = \hat{\psi}(1) = 0, \quad \hat{\psi}(\zeta) > 0 \text{ for } \zeta \neq 0, 1.$$  

The precise assumptions on $\hat{\psi}$, which we need in our investigations, are stated in Theorem 2.3. This completes the formulation of the model. 

(1.1) and (1.2) are the equations of linear elasticity theory. This subsystem is coupled to the Allen-Cahn equation (1.3), which governs the evolution of the order parameter $S$. 

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The system \((1.1) - (1.3)\) satisfies the second law of thermodynamics. More precisely, the Clausius-Duhem inequality is satisfied with the free energy

\[ \psi^{*}(\varepsilon, S) = \mathcal{W}(\varepsilon, S) + \frac{1}{\mu^{1/2}} \hat{\psi}(S) + \frac{\mu^{1/2} \lambda}{2} |\nabla x S|^2. \]  

(1.9)

From this expression we see that the parameter \(\lambda\) determines the energy density of the phase interface. We assume that the parameters \(\mu\) and \(\lambda\) vary in intervals \((0, \mu_0]\) and \((0, \lambda_0]\), respectively, with \(\mu_0 > 0\) and \(\lambda_0 > 0\) chosen sufficiently small. The scaling \(\frac{\varepsilon}{(\mu \lambda)^{1/2}}\) on the right hand side of \((1.3)\) is necessary for otherwise the propagation speed of the diffuse interface would tend to zero for \(\mu \to 0\) or \(\lambda \to 0\).

We give now a slightly sketchy overview of the main results in this article. The precise definitions and statements are given in Section 2 and in particular in Section 2.3.

We denote solutions of the Allen-Cahn equation by \((u_{AC}, T_{AC}, S_{AC})\), to distinguish them from approximate solutions, which we construct later. The values \(S_{AC}(t, x) \approx 0\) or \(S_{AC}(t, x) \approx 1\) indicate that at the point \(x \in \Omega\) at time \(t\) the crystal structure of the material of the solid belongs to phase 1 or to phase 2, respectively. The set of all \(x \in \Omega\) with \(0 < S_{AC}(t, x) < 1\) is the region of the diffuse interface at time \(t\). The level set

\[ \Gamma_{AC}(t) = \{ x \in \Omega \mid S_{AC}(t, x) = \frac{1}{2} \} \]  

(1.10)

belongs to this region. For \(x \in \Gamma_{AC}(t)\) we denote by \(s_{AC}(t, x) \in \mathbb{R}\) the normal speed of this level set at \(x\), and we call \(s_{AC}\) the speed of the diffuse interface. For this speed we derive an expression of the form

\[ s_{AC}(t, x) = (s_{00} + \lambda^{1/2} s_{01}) + \mu^{1/2} (s_{10} + \lambda^{1/2} s_{11}) + |\ln \mu|^{3} \mu s^{(\mu \lambda)}(\varepsilon, S). \]  

(1.11)

We call this expression the kinetic relation of the Allen-Cahn model. It is the central result of this paper. The remainder term \(s^{(\mu \lambda)}(\varepsilon, S)\) depends on \(\mu\) and \(\lambda\), but in Section 6 we prove that there exist numbers \(\mu_0 > 0\) and \(\lambda_0 > 0\) and a constant \(C_\varepsilon\) such that the \(L^2\)-norm satisfies

\[ \|s^{(\mu \lambda)}(\varepsilon, S)\|_{L^2_{\Gamma_{AC}(t)}} \leq C_\varepsilon \]  

(1.12)

for all \(0 < \mu \leq \mu_0\) and all \(0 < \lambda \leq \lambda_0\), where \(s^{(\mu \lambda)}(\varepsilon, S)\) denotes the function \(x \mapsto s^{(\mu \lambda)}(\varepsilon, S, x) : \Gamma_{AC}(t) \to \mathbb{R}\). Therefore \((1.11)\) is an asymptotic expansion for the propagation speed of the diffuse interface, which is uniformly valid with respect to the parameters \(\mu\) and \(\lambda\). For sufficiently small \(\mu\) the leading term \(s_0(t, x) = s_{00} + \lambda^{1/2} s_{01}\) and the second term \(\mu^{1/2} s_1(t, x) = \mu^{1/2} (s_{10} + \lambda^{1/2} s_{11})\) dominate over the remainder term \(|\ln \mu|^{3} \mu s^{(\mu \lambda)}(\varepsilon, S)\). We can therefore read off the behavior of the Allen-Cahn model with respect to the parameters \(\mu\) and \(\lambda\) from the first two terms in \((1.11)\).

The terms \(s_{00} + \lambda^{1/2} s_{01}\) and \(s_{10} + \lambda^{1/2} s_{11}\) are explicitly given in Theorem 2.3. We restrict ourselves here to state the form of the leading term. We assume that \(0 \leq t_1 < t_2 < \infty\) are given times. We study the propagation of the interface for \(t\) varying in the interval \([t_1, t_2]\). For \(t\) from this interval the leading term is

\[ s_0(t, x) = s_{00} + \lambda^{1/2} s_{01} = \frac{c}{c_1} \left( -\sigma : \langle \hat{T} \rangle(t, x) + \lambda^{1/2} c_1 \kappa_\Gamma(t, x) \right), \]  

(1.13)
where \( c \) is the mobility constant from (1.3), 
\[
c_1 = \int_0^1 \sqrt{2 \hat{\psi}(\vartheta)} \, d\vartheta
\]
is computed from the double well potential, \( \kappa_{\Gamma}(t, x) \) is twice the mean curvature of the surface \( \Gamma_{AC}(t) \) at the point \( x \), and \( \hat{T} \) is the stress field in the solution \( (t, x) \mapsto (\hat{u}(t, x), \hat{T}(t, x)) \) of the transmission problem:

\[
\begin{align*}
-\text{div}_x \hat{T} &= b, \quad \text{(1.14)} \\
\hat{T} &= D(\varepsilon(\nabla_x \hat{u}) - \varepsilon \hat{S}), \quad \text{(1.15)} \\
[\hat{u}] &= 0, \quad \text{(1.16)} \\
[\hat{T}]n &= 0, \quad \text{(1.17)} \\
\hat{u}(t)|_{\partial \Omega} &= U(t), \quad \text{(1.18)}
\end{align*}
\]

Here \( (t, x) \mapsto \hat{S}(t, x) \) is a function, which takes only the values 0 or 1 and jumps across the interface \( \Gamma_{AC} = \{(t, x) \in [t_1, t_2] \times \Omega \mid x \in \Gamma_{AC}(t)\} \).

\([\hat{u}]\) and \([\hat{T}]\) are the jumps of the functions \( \hat{u} \) and \( \hat{T} \) across \( \Gamma_{AC} \), and \( b, U \) are the volume force and boundary data from equations (1.1) and (1.4). The equations (1.14), (1.15) hold in the domain \( ([t_1, t_2] \times \Omega) \setminus \Gamma_{AC} \). For every fixed \( t \) the problem (1.14) – (1.18) is an elliptic transmission problem for the function \( x \mapsto (\hat{u}(t, x), \hat{T}(t, x)) \) in the domain \( \Omega \). Therefore \( t \) can be considered to be a parameter in this problem. Finally, the expression

\[
\langle \hat{T} \rangle = \frac{1}{2} (\hat{T}^{(+)}) + \hat{T}^{(-)}
\]
in equation (1.13) denotes the mean value of the values \( \hat{T}^{(+)} \) and \( \hat{T}^{(-)} \) on both sides of the interface \( \Gamma_{AC} \).

The terms \( s_{10} \) and \( s_{11} \) are determined by a more complicated transmission problem, for which the transmission conditions are also posed on the interface \( \Gamma_{AC} \), and \( b, U \) are the volume force and boundary data from equations (1.1) and (1.4). The equations (1.14), (1.15) hold in the domain \( ([t_1, t_2] \times \Omega) \setminus \Gamma_{AC} \). For every fixed \( t \) the problem (1.14) – (1.18) is an elliptic transmission problem for the function \( x \mapsto (\hat{u}(t, x), \hat{T}(t, x)) \) in the domain \( \Omega \). Therefore \( t \) can be considered to be a parameter in this problem. Finally, the expression

\[
\langle \hat{T} \rangle = \frac{1}{2} (\hat{T}^{(+)}) + \hat{T}^{(-)}
\]
in equation (1.13) denotes the mean value of the values \( \hat{T}^{(+)} \) and \( \hat{T}^{(-)} \) on both sides of the interface \( \Gamma_{AC} \).

The terms \( s_{10} \) and \( s_{11} \) are determined by a more complicated transmission problem, for which the transmission conditions are also posed on the interface \( \Gamma_{AC} \), and by a coupled system of ordinary differential equations for two functions \( S_0 \) and \( S_1 \), which are needed in the construction of an asymptotic solution of the Allen-Cahn model (1.1) – (1.5). The coefficients of the second transmission problem and of the system of ordinary differential equations depend on the solution \( (\hat{u}, \hat{T}) \) of (1.14) – (1.18). Both transmission problems together can thus be considered to be a larger transmission problem, which is recursively solvable.

The derivation of the kinetic relation (1.11) is based on the construction of an asymptotic solution \( (u^{(\mu\lambda)}, T^{(\mu\lambda)}, S^{(\mu\lambda)}) \) for the Allen-Cahn model. From this asymptotic solution it is seen that the width of the diffuse interface in the Allen-Cahn model is proportional to the parameter

\[
B = (\mu\lambda)^{1/2}.
\]

(1.19)

We call \( B \) the interface width parameter. As will be explained in Section 2.4, the interface energy density is proportional to the parameter

\[
E = \lambda^{1/2}.
\]

(1.20)

We call \( E \) the interface energy parameter.

The kinetic relation (1.11) and the equation (1.19) together have consequences for the efficiency of numerical simulations of interfaces with low interface energy density, which we sketch here. A precise discussion is given in Section 2.4.
The explicit expressions in Theorem 2.3 show that the second term \( \mu^{1/2}(s_{10} + \lambda^{1/2}s_{11}) \) in (1.11) is of a very special form. We therefore argue that this term does not have a physical meaning, only the leading term \( s_{00} + \lambda^{1/2}s_{01} \) is physically relevant. This means that in (1.11) the term

\[
\mathcal{E}^{(\mu\lambda)} = \mu^{1/2}(s_{10} + \lambda^{1/2}s_{11}) + |\ln \mu|^{3} \mu s_{\infty}^{(\mu\lambda)}
\]

is a mathematical error term, which in a precise numerical simulation of the evolution of the interface must be made small by choosing \( \mu^{1/2} \) small enough. Therefore we call \( \mathcal{E}^{(\mu\lambda)} \) the model error and \( F = \mu^{1/2} \) the error parameter. By (1.19), the interface width is proportional to the error parameter \( F \). This means that the interface width depends on the size of the model error. The smaller the model error is, which we want to allow, the smaller the interface width must be chosen. The total error in a numerical simulation consists of the model error and the numerical error. In order to make the numerical error small, the grid spacing must be chosen small enough to resolve the transition of the order parameter from 0 to 1 across the diffuse interface. When the interface width is small, we must therefore choose the grid spacing small, which means that the numerical effort is high.

From (1.19) and (1.20) we see that for constant values of the error parameter \( F \) the interface width is proportional to the interface energy density parameter \( E \). Thus, when we want to precisely simulate an interface with small interface energy density, we must choose small values for \( E \) and \( F \), hence the interface width \( \mathcal{F} = EF \) becomes very small. As a consequence, also the grid spacing must be chosen very small, which means that numerical simulations of interfaces with low interface energy based on the Allen-Cahn model are not efficient. Of course, the efficiency can be improved by using adaptive mesh refinement and other advanced numerical techniques, but still it would be advantageous if such tools could be avoided.

Often the Allen-Cahn model is formulated using the parameters \( E \) and \( F \) instead of \( \mu \) and \( \lambda \). It might therefore be helpful to shortly discuss the form, which our results take when this formulation is used. With these parameters the Allen-Cahn equation (1.3) is

\[
\partial_t S = -\frac{c}{B} \left( \partial_S \mathcal{W}(\varepsilon(\nabla_x u), S) + E \left( \frac{\dot{\psi}(S)}{B} - B \Delta_x S \right) \right),
\]

the free energy (1.9) becomes

\[
\psi_{EB}(\varepsilon, S) = \mathcal{W}(\varepsilon, S) + E \left( \frac{\dot{\psi}(S)}{B} + \frac{B}{2} |\nabla_x S|^2 \right),
\]

and the kinetic relation (1.11) takes the form

\[
s_{AC}(t, x) = (s_{00} + E s_{01}) + \frac{B}{E} (s_{10} + E s_{11}) + \left| 2 \ln \left( \frac{B}{E} \right) \right|^{3} \left( \frac{B}{E} \right)^{2} s_{\infty}^{(BE)}.
\]

From this equation we see that the model error

\[
\mathcal{E}^{(BE)} = \frac{B}{E} (s_{10} + E s_{11}) + \left| 2 \ln \left( \frac{B}{E} \right) \right|^{3} \left( \frac{B}{E} \right)^{2} s_{\infty}^{(BE)}
\]

is governed by the ratio \( \frac{B}{E} \). When one reduces the interface energy density \( E \) and one wants to keep the model error constant, one must reduce the interface width \( B \) by the
same proportion. In a simulation of an interface without interface energy, the total model error is the sum \( E_{\text{total}}^{(BE)} = E_{s_{01}} + E_{(BE)} \). To reduce the error in such a simulation we must reduce \( E \) and the fraction \( \frac{B}{E} \), hence \( B \) must be reduced faster than \( E \).

The paper is organized as follows. The main results are presented in Section 2, where we first state the transmission problems and the system of ordinary differential equations, whose solutions are needed to compute the coefficients \( s_{00}, \ldots, s_{11} \) in (1.11). These coefficients are explicitly given in Theorem 2.3. Moreover, in this theorem we also state properties of the asymptotic solution \( (u(\mu\lambda), T(\mu\lambda), S(\mu\lambda)) \), which is constructed in later sections. In particular, we state the scaling law (1.19) for the width of the diffuse interface. These properties are needed in Section 2.4 where we precisely discuss the model error and the numerical efficiency. The estimate (1.12), which is the most important mathematical result of this paper, is stated in Theorem 2.8.

Sections 3 – 5 contain the proof of Theorem 2.3. In Section 3 we construct the approximate solution \( (u(\mu\lambda), T(\mu\lambda), S(\mu\lambda)) \). That is, we state the inner and outer expansions which define the function \( (u(\mu\lambda), T(\mu\lambda), S(\mu\lambda)) \). In these asymptotic expansions functions appear, which are obtained as solutions of systems of algebraic and differential equations. These systems are also stated in Section 3. The system for the outer expansion can be readily solved, and the solution of the system of ordinary differential equations for the inner expansion is more involved and is discussed in Section 4. In two equations of this system a linear differential operator appears with kernel different from \( \{0\} \). In order that these differential equations be solvable the right hand sides must satisfy orthogonality conditions. The right hand sides contain the coefficients \( s_{00}, \ldots, s_{11} \) of the kinetic relation (1.11). The orthogonality conditions dictate the form of these coefficients; therefrom the equation (1.11) results. In Section 3 we verify that \( (u(\mu\lambda), T(\mu\lambda), S(\mu\lambda)) \) is really an asymptotic solution of the model equations (1.1) – (1.3) and prove the necessary estimates.

In Section 6 we prove the estimate (1.12). The proof uses the residue, with which the function \( (u(\mu\lambda), T(\mu\lambda), S(\mu\lambda)) \) satisfies the equations (1.1) – (1.3). The main difficulty in the proof is that though we want to prove that \( s_{\mu\lambda}^{\infty} \) is bounded uniformly with respect to \( \mu \) and \( \lambda \), the residue term itself is not bounded for \( \lambda \to 0 \), but instead behaves like \( \lambda^{-1/2} \).

In the bibliography of [8] we gave many references to the literature on existence, uniqueness and asymptotics for models containing the Allen-Cahn and Cahn-Hilliard equations. We refer the reader to that bibliography and discuss here only some publications, which are of interest in the construction of asymptotic solutions.

We believe that for the model (1.1) – (1.3) an asymptotic solution was constructed and used to identify the associated sharp interface problem for the first time in [23], following earlier such investigations for other phase field models. For example, in [13] these investigations were carried out for a model from solidification theory, which consists of the Allen-Cahn equation coupled to the heat equation.

The considerations in [13, 23] are formal, since it is not shown that the asymptotic solution converges to an exact solution of the model equations for \( \mu \to 0 \). Under the assumption that the associated sharp interface problems have smooth solutions, this was proved in [28] for the Allen-Cahn equation, in [10] for the Cahn-Hilliard equation, and in [11] for the model from solidification theory and in [14] for a model consisting of the Cahn-Hilliard equation coupled with the elasticity equations. The proofs use variants of a spectral estimate derived in [16]. For the model from solidification theory the associated sharp interface model is the Mullins-Sekerka model with surface tension.
In [15] an asymptotic solution for the Cahn-Hilliard equation has recently been constructed with a method different from the one used in [10], and which is similar to our method.

In [25] the numerical efficiency of simulations based on the phase field model consisting of the Allen-Cahn equation coupled to the heat equation is studied. It is shown that for suitably chosen coefficients of the model the second order term in an asymptotic expansion of the solution vanishes. By arguments similar to the ones we gave in the above discussion it is seen that this improves the numerical efficiency of the model. This result has been improved and generalized in [9, 17, 24]. A similar idea is also present in [22].

Since the construction of asymptotic solutions is based on sharp interface problems, a rigorous analysis of these problems is of special interest. Of particular interest is the Hele-Shaw problem with surface tension, since this is the sharp interface problem associated with the Cahn-Hilliard equation. Existence, uniqueness, and regularity of classical solutions of this problem have been investigated in [18, 19, 20]. In [21] it is shown that if the initial data are close to a sphere then a classical solution exists and converges to spheres. Existence of solutions to the Mullins-Sekerka problem mentioned above has been shown in [27].

The model (1.1) – (1.6) describes the evolution of phase transitions in a solid when temperature effects are negligible. This model is the prototype of a large class of models obtained by extensions and generalizations of the model, which are used in the engineering sciences to simulate the behavior of complex and functional materials. From the very large literature in this field we cite here only [11, 29, 30, 31, 32].

2 The kinetic relation

2.1 Notations

For given fixed times \( 0 \leq t_1 < t_2 < \infty \) let

\[ Q = [t_1, t_2] \times \Omega \subseteq \mathbb{R}^4. \]

The construction of the asymptotic solution \((u(\mu), T(\mu), S(\mu))\) is based on a surface \(\Gamma(\mu)(t)\), which for \( t_1 \leq t \leq t_2 \) moves in \(\Omega\) and which will be the level set

\[ \Gamma(\mu)(t) = \{ x \in \Omega \mid S(\mu)(t, x) = \frac{1}{2} \}. \]

We set

\[ \Gamma(\mu) = \{(t, x) \in Q \mid x \in \Gamma(\mu)(t)\}. \tag{2.1} \]

To simplify the notation we often drop the index \( \lambda \) or both indices \( \mu \) and \( \lambda \) and write \(\Gamma(\mu)\) and \(\Gamma(\mu)\) or simply \(\Gamma\) and \(\Gamma\). Similarly, we often write \((u(\mu), T(\mu), S(\mu))\) for the asymptotic solution and use the same convention also in other notations. Both indices are specified if the dependence on \( \lambda \) becomes important.

The precise definition of the family \(\{\Gamma(\mu)(t)\}_{t_1 \leq t \leq t_2}\) is given in the next section, and in Section 2.4 we associate \(\Gamma(\mu)(t)\) with the level set \(\Gamma_{\text{AC}}(\mu)(t) = \Gamma_{\text{AC}}(t)\) introduced in (1.10). To introduce notations we assume here that \(\Gamma\) is a known, orientable, three dimensional \(C^k\)-manifold with \(k \geq 1\) sufficiently large embedded in \(Q\) such that \(\Gamma(t)\) is a regular two dimensional surface in \(\Omega\) for every \( t \in [t_1, t_2] \). Let

\[ n : \Gamma \to \mathbb{R}^3 \]  \tag{2.2}
be a continuous vector field such that \( n(t, x) \in \mathbb{R}^3 \) is a unit normal vector to \( \Gamma(t) \) at \( x \in \Gamma(t) \), for every \( t \in [t_1, t_2] \). For \( \delta > 0 \) and \( t \in [t_1, t_2] \) define the sets

\[
U_\delta(t) = \{ x \in \Omega \mid \text{dist}(x, \Gamma(t)) < \delta \} \quad \text{and} \quad \mathcal{U}_\delta = \{ (t, x) \in Q \mid x \in U_\delta(t) \}. 
\]  

(2.3)

We assume that there is \( \delta > 0 \) such that \( \mathcal{U}_\delta \subseteq Q \). Since \( \Gamma \) is a regular \( C^1 \)-manifold in \( Q \), then \( \delta \) can be chosen sufficiently small such that for all \( t \in [t_1, t_2] \) the mapping \( (\eta, \xi) \mapsto x(t, \eta, \xi) = \eta + \xi n(t, \eta) : \Gamma(t) \times (-\delta, \delta) \to \mathcal{U}_\delta(t) \) \( \text{is bijective. \ We say that this mapping defines new coordinates \((\eta, \xi)\) in } \mathcal{U}_\delta(t) \text{ and } (t, \eta, \xi) \text{ in } \mathcal{U}_\delta. \) If no confusion is possible we switch freely between the coordinates \((t, x)\) and \((t, \eta, \xi)\). In particular, if \( (t, x) \mapsto w(t, x) \) is a function defined on \( \mathcal{U}_\delta \) we write \( w(t, \eta, \xi) \) for \( w(t, x(t, \eta, \xi)) \), as usual.

We use the standard convention and denote for a function \( w \) defined on a subset \( U \) of \( Q \) by \( w(t) \) the function \( x \mapsto w(t, x) \), which is defined on the set \( \{ x \mid (t, x) \in U \} \subseteq \mathbb{R}^3 \).

If \( w \) is a function defined on \( \mathcal{U}_\delta(t) \setminus \Gamma(t) \), we set for \( \eta \in \Gamma(t) \)

\[
w^{(\pm)}(\eta) = \lim_{\xi \to 0} w(\eta \pm \xi n(t, \eta)), \\
(\partial_n^i w)^{(+)}(\eta) = \lim_{\xi \to 0} \frac{\partial}{\partial \xi^i} w(\eta + \xi n(t, \eta)), \quad i \in \mathbb{N}, \\
(\partial_n^i w)^{(-)}(\eta) = \lim_{\xi \to 0} \frac{\partial}{\partial \xi^i} w(\eta + \xi n(t, \eta)), \quad i \in \mathbb{N}, \\
[w](\eta) = w^{(0)}(\eta) - w^{(-)}(\eta), \\
[\partial_n^i w](\eta) = (\partial_n^i w)^{(+))(\eta)} - (\partial_n^i w)^{(-)}(\eta), \\
\langle w \rangle(\eta) = \frac{1}{2}(w^{(+))(\eta)} + w^{(-)}(\eta)),
\]

provided that the one-sided limits in these equations exist. If \( w \) is defined on \( \mathcal{U}_\delta \setminus \Gamma \), we set

\[
w^{(\pm)}(t, \eta) = \left( w^{(\pm)}(t) \right)(\eta), \quad (\partial_n^i w)^{(\pm)}(t, \eta) = \left( (\partial_n^i w)^{\pm}(t) \right)(\eta),
\]

and define \([w](t, \eta), \langle w \rangle(t, \eta), [\partial_n^i w](t, \eta)\) as above. Let \( \tau_1(\eta), \tau_2(\eta) \in \mathbb{R}^3 \) be two orthogonal unit vectors to \( \Gamma(t) \) at \( \eta \in \Gamma(t) \). For functions \( w : \Gamma(t) \to \mathbb{R}, W : \Gamma(t) \to \mathbb{R}^3 \) we define the surface gradients by

\[
\nabla_\Gamma w = (\partial_{\tau_1} w)\tau_1 + (\partial_{\tau_2} w)\tau_2, \\
\nabla_\Gamma W = (\partial_{\tau_1} W) \otimes \tau_1 + (\partial_{\tau_2} W) \otimes \tau_2,
\]

(2.5) \( (2.6) \)

where for vectors \( c, d \in \mathbb{R}^3 \) a \( 3 \times 3 \)-matrix is defined by

\[
c \otimes d = (c_i d_j)_{i,j=1,2,3}.
\]

With \( (2.5), (2.6) \) we have for functions \( w : \mathcal{U}_\delta(t) \to \mathbb{R} \) and \( W : \mathcal{U}_\delta(t) \to \mathbb{R}^3 \) at \( \eta \in \Gamma(t) \) the decompositions

\[
\nabla_x w = (\partial_n w)n + \nabla_\Gamma w, \\
\nabla_x W = (\partial_n W) \otimes n + \nabla_\Gamma W,
\]

(2.7) \( (2.8) \)
where \( n = n(t, \eta) \) is the unit normal vector to \( \Gamma(t) \).

The normal speed of the family of surfaces \( t \mapsto \Gamma(t) \) is of fundamental importance in this paper. Therefore we give a precise definition.

**Definition 2.1** Let \( m(t, \eta) = (m'(t, \eta), m''(t, \eta)) \in \mathbb{R} \times \mathbb{R}^3 \) be a normal vector to \( \Gamma \) at \( (t, \eta) \in \Gamma \). The normal speed of the family of surfaces \( t \mapsto \Gamma(t) \) at \( \eta \in \Gamma(t) \) is defined by

\[
s(t, \eta) = \frac{-m'(t, \eta)}{m''(t, \eta) \cdot n(t, \eta)},
\]

with the unit normal vector \( n(t, \eta) \in \mathbb{R}^3 \) to \( \Gamma(t) \).

Note that with this definition the speed is measured positive in the direction of the normal vector field \( n \). Since \( m''(t, \eta) \in \mathbb{R}^3 \) is a normal vector to \( \Gamma(t) \), the denominator in (2.9) is different from zero.

If \( \omega = (\omega', \omega'') \in \mathbb{R} \times \mathbb{R}^3 \) is a tangential vector to \( \Gamma \) at \( (t, \eta) \) with \( \omega' \neq 0 \), then with the unit normal \( n(t, \eta) \in \mathbb{R}^3 \) to \( \Gamma(t) \) the vector \( -\omega'' \cdot n, \omega' n \) is a normal vector to \( \Gamma \) at \( (t, \eta) \), hence (2.9) implies that the normal speed at \( \eta \in \Gamma(t) \) is given by

\[
s(t, \eta) = \frac{n \cdot \omega''}{\omega' n \cdot n} = \frac{n \cdot \omega''}{\omega'}.
\]

(2.10)

For later use we prove the following

**Lemma 2.2** Let \( x \in \mathcal{U}_\delta(t_0) \) be a point having the representation \( x = \eta + n(t, \eta)\xi \) in the \((\eta, \xi)-\)coordinates, where \( \eta = \eta(t, x) \in \Gamma(t) \) and \( \xi = \xi(t, x) \). Then the normal speed satisfies

\[
s(t_0, \eta) = n(t_0, \eta) \cdot \partial_t \eta(t_0, x) = -\partial_t \xi(t_0, x).
\]

(2.11)

The tangential component of the vector \( \partial_t \eta(t_0, x) \in \mathbb{R}^3 \) to the surface \( \Gamma(t_0) \) is equal to \(-\xi \partial_t n(\eta(t_0, x), \eta(t_0, x)).\)

**Proof:** By definition of \( \mathcal{U}_\delta \), there is a neighborhood \( U \) of \( t_0 \) in \([t_1, t_2]\) such that \( \{x\} \times U \subseteq \mathcal{U}_{\delta} \), which implies that \( x \) has the representation

\[
x = \eta(t, x) + \xi(t, x)n(t, \eta(t, x)).
\]

for all \( t \in U \). We differentiate this equation and obtain

\[
0 = \partial_t x = n \partial_t \xi + \xi \partial_t n + \partial_t \eta.
\]

(2.12)

From \( 0 = \partial_t 1 = \partial_t|n|^2 = 2n \cdot \partial_t n \) we see that \( \partial_t n \) is tangential to \( \Gamma(t) \), hence (2.12) implies that the tangential component of \( \partial_t \eta \) is equal to \(-\xi \partial_t n \). Multiplication of (2.12) with \( n \) yields

\[
\partial_t \xi = -n \cdot \partial_t \eta.
\]

(2.13)

Since \( \partial_t(t, \eta(t, x)) = (1, \partial_t \eta(t, x)) \) is a tangential vector to \( \Gamma \), it follows from (2.10) that \( s = \frac{n \cdot \partial_t n}{\omega'} = n \cdot \partial_t \eta \), which together with (2.13) implies (2.11). \( \blacksquare \)
2.2 The evolution problem for the level set $\Gamma(\mu)$

The level set $\Gamma = \Gamma(\mu)$ of $S(\mu)$ defined in (2.1) is determined by an evolution problem for the family of surfaces $t \mapsto \Gamma(t)$. To state this evolution problem let $\mathcal{N}$ be the operator, which assigns the normal speed to the family $t \mapsto \Gamma(t)$, i.e.

$$ s(t, x) = \mathcal{N}(\Gamma)(t, x), $$

with $s(t, x) = s(\mu)(t, x)$ defined by (2.9). The evolution problem is given by

$$ \mathcal{N}(\Gamma)(t) = \mathcal{K}(\mu)(\Gamma(t)), \quad t_1 \leq t \leq t_2, \quad (2.14) $$

where $\mathcal{K}(\mu)$ is the non-local evolution operator, which has the form

$$ \mathcal{K}(\mu)(\Gamma(t))(x) = s_0(\hat{T}, \kappa_1, \lambda^{1/2})(t, x) + \mu^{1/2}s_1(\hat{u}, \hat{T}, T, S_0, S_1, \lambda^{1/2})(t, x), \quad (2.15) $$

for $x \in \Gamma(t)$. Here $(\hat{u}, \hat{T}, \hat{u}, \hat{T}, S_0, S_1)$ is the solution of a transmission-boundary value problem for a coupled system of elliptic partial differential equations and ordinary differential equations, which can be solved recursively. $\kappa_1(t, x)$ denotes twice the mean curvature of the surface $\Gamma(t)$ at $x \in \Gamma(t)$. With the principle curvatures $\kappa_1, \kappa_2$ of $\Gamma(t)$ at $x \in \Gamma(t)$ we thus have

$$ \kappa_1(t, x) = \kappa_1(t, x) + \kappa_2(t, x). $$

The transmission condition is posed on $\Gamma(t)$. Therefore the functions $\hat{u}, \hat{T}, \hat{u}, \hat{T}$ and $S_1$ depend on $\Gamma(t)$. We first state and discuss the transmission-boundary value problem. The precise form of the functions $s_0$ and $s_1$ is given in Theorem 2.3 following below.

Let $\hat{S} : Q \setminus \Gamma \rightarrow \{0, 1\}$ be a piecewise constant function, which only takes the values 0 and 1 with a jump across $\Gamma$. The sets

$$ \gamma = \{(t, x) \in Q \mid \hat{S}(t, x) = 0\}, \quad \gamma(t) = \{x \in \Omega \mid \Gamma(t) \mid (t, x) \in \gamma\}, $$

$$ \gamma' = \{(t, x) \in Q \mid \hat{S}(t, x) = 1\}, \quad \gamma'(t) = \{x \in \Omega \mid \Gamma(t) \mid (t, x) \in \gamma'\} $$

yield partitions $Q = \gamma \cup \Gamma \cup \gamma'$ and $\Omega = \gamma(t) \cup \Gamma(t) \cup \gamma'(t)$ of $Q$ and $\Omega$, respectively. If $x$ belongs to $\gamma(t)$ or $\gamma'(t)$, then the crystal structure at the material point $x$ at time $t$ belongs to phase 1 or phase 2, respectively. We assume that the normal vector field $n$ given in (2.2) is such that the vector $n(t, x)$ points into the set $\gamma'(t)$ for every $x \in \Gamma(t)$.

The transmission-boundary value problem can be separated into two transmission-boundary value problems for the elasticity equations and a boundary value problem for a coupled system of two ordinary differential equations. To state the complete problem we fix $t \in [t_1, t_2]$ and assume that $\Gamma(t)$ is known. In the first transmission-boundary problem the unknowns are the displacement $x \mapsto \hat{u}(t, x) \in \mathbb{R}^3$ and the stress tensor $x \mapsto \hat{T}(t, x) \in \mathcal{S}^3$, which must satisfy the equations

$$ -\text{div}_x \hat{T} = b, \quad (2.16) $$

$$ \hat{T} = D(\varepsilon(\nabla_x \hat{u}) - \mathcal{I} \hat{S}), \quad (2.17) $$

$$ [\hat{u}] = 0, \quad (2.18) $$

$$ [\hat{T}]n = 0, \quad (2.19) $$

$$ \hat{u}(t)|_{\partial \Omega} = U(t), \quad (2.20) $$

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with \( b \) and \( U \) given in (1.1) and (1.4). In the second transmission-boundary problem the unknowns are the displacement \( x \mapsto \hat{u}(t, x) \in \mathbb{R}^3 \) and the stress tensor \( x \mapsto \hat{T}(t, x) \in S^3 \), and the problem is

\[
-\text{div}_x \hat{T} = 0, \quad (2.21)
\]

\[
\hat{T} = D \left( \varepsilon(\nabla_x \hat{u}) - \varepsilon \frac{\hat{T}}{\psi''(S)} \right), \quad (2.22)
\]

\[
[\hat{u}] = 0, \quad (2.23)
\]

\[
[\hat{T}]_n = 0, \quad (2.24)
\]

\[
\hat{u}(t)|_{\partial\Omega} = 0. \quad (2.25)
\]

The equations (2.16), (2.17) and (2.21), (2.22) must hold on the set \( \Omega \setminus \Gamma(t) \), whereas the equations (2.18), (2.19) and (2.23), (2.24) are posed on \( \Gamma(t) \).

In the boundary value problem for the ordinary differential equations the unknowns are \( S_0 : \mathbb{R} \to \mathbb{R} \), \( S_1 : \Gamma \times \mathbb{R} \to \mathbb{R} \) and \( s_0 : \Gamma \to \mathbb{R} \). We use the notations \( S'_0(t, \eta, \zeta) = \partial_\zeta S_1(t, \eta, \zeta) \), \( S''_0(t, \eta, \zeta) = \partial^2_\zeta S'_0(t, \eta, \zeta) \). In this problem not only \( t \), but also \( \eta \in \Gamma(t) \) is a parameter. For all \( \zeta \in \mathbb{R} \) and all values of the parameter \( \eta \in \Gamma(t) \) the unknowns must satisfy the coupled ordinary differential equations

\[
\psi'(S_0(\zeta)) - S''_0(\zeta) = 0, \quad (2.26)
\]

\[
\psi''(S_0(\zeta)) S_1(t, \eta, \zeta) - S''_1(t, \eta, \zeta) = F_1(t, \eta, \zeta), \quad (2.27)
\]

and the boundary conditions

\[
S_0(0) = \frac{1}{2}, \quad \lim_{\zeta \to -\infty} S_0(\zeta) = 0, \quad \lim_{\zeta \to \infty} S_0(\zeta) = 1, \quad (2.28)
\]

\[
\lim_{\zeta \to -\infty} S_1(t, \eta, \zeta) = \frac{\varepsilon : \hat{T}(-)(t, \eta)}{\psi''(0)}, \quad (2.29)
\]

\[
\lim_{\zeta \to +\infty} S_1(t, \eta, \zeta) = \frac{\varepsilon : \hat{T}(+)(t, \eta)}{\psi''(1)}, \quad (2.30)
\]

\[
S_1(t, \eta, 0) = 0, \quad (2.31)
\]

with the right hand side of (2.27) given by

\[
F_1(t, \eta, \zeta) = \varepsilon : \left( [\hat{T]}(t, \eta) S_0(\zeta) + \hat{T}(-)(t, \eta) \right)
\]

\[
+ \left( \frac{s_0(t, \eta)}{c} - \lambda^{1/2} \kappa \Gamma(t, \eta) \right) S'_0(\zeta), \quad (2.32)
\]

where the constant \( c > 0 \) is given in (1.3).

The linear elliptic system (2.16), (2.17) differs from the standard elasticity system only by the term \(-D\varepsilon S\). This term is known since \( \Gamma(t) \) is given. Under suitable regularity assumptions for the given functions \( b \) and \( U \) and very mild assumptions on the regularity of the interface \( \Gamma(t) \) the problem has a unique weak solution \((\hat{u}, \hat{T})\). This can be proved by standard methods from functional analysis. Of course, the regularity of the solution depends on the regularity of \( b \), \( U \) and \( \Gamma(t) \).
After insertion of the stress tensor $\hat{T}$ from this solution into (2.22), the equations (2.21) – (2.25) form a transmission-boundary value problem of the same type as (2.16) – (2.20), with unique solution $(\hat{u}, \hat{T})$ determined by the same methods.

We also insert $\hat{T}$ into (2.29), (2.30) and (2.32), which determines the right hand side of the differential equation (2.27) and the boundary conditions (2.29), (2.30) posed at $\pm \infty$. The nonlinear differential equation (2.26) has a unique solution $S_0$ satisfying the boundary conditions (2.28). By insertion of $S_0$ into (2.27) and (2.32), equation (2.27) becomes a linear differential equation for $S_1$, however with an additional unknown function $s_0$ in the right hand side. This function is constant with respect to $\zeta$. We sketch here the procedure used to determine $s_0$. This procedure is standard in investigations of the asymptotics of phase field models:

The second order differential operator $(\hat{\psi}''(S_0) - \partial_\zeta^2)w = F_1$ has a solution $w \in L^2(\mathbb{R})$ if and only if the orthogonality condition

$$\int_{-\infty}^{\infty} F_1(t, \eta, \zeta) S_0'(\zeta) d\zeta = 0 \quad (2.33)$$

holds. It turns out that though the function $F_1$ defined in (2.32) does not in general belong to $L^2(\mathbb{R})$ and the solution $S_1(t, \eta, \cdot)$ is not sought in $L^2(\mathbb{R})$, which is seen by differentiating the equation (2.26). From functional analysis we thus know that for $F_1 \in L^2(\mathbb{R})$ the differential equation $(\hat{\psi}''(S_0) - \partial_\zeta^2)w = F_1$ a solution $w \in L^2(\mathbb{R})$ if and only if the orthogonality condition (2.33) is sufficient for the solution $S_1$ to exist. Comparison with (2.32) shows that (2.33) can be satisfied by choosing the constant $s_0(t, \eta)$ suitably. This defines the function $s_0 : \Gamma \to \mathbb{R}$ uniquely. Since $F_1$ depends on $\hat{T}$, $\kappa_\Gamma$ and $\lambda^{1/2}$, it follows that also $s_0$ is a function of these variables:

$$s_0(t, \eta) = s_0(\hat{T}, \kappa_\Gamma, \lambda^{1/2})(t, x).$$

The explicit expression for $s_0$ obtained in this way is stated below in (2.40). In fact, $s_0(\hat{T}, \kappa_\Gamma, \lambda^{1/2})$ is the first term on the right hand side in the expression (2.15) for $K^{(\mu)}(\Gamma(t))$.

The procedure sketched here is discussed precisely in Section 4.2 when we determine the second term $s_1$ in (2.15), which is obtained from a similar, but more complicated boundary value problem.

2.3 The asymptotic solution and the kinetic relation

To state the properties of the asymptotic solution and the kinetic relation in Theorem 2.3 we introduce some definitions.

We need in our investigations that the second derivatives $\hat{\psi}''(0)$ and $\hat{\psi}''(1)$ of the double well potential at the minima 0 and 1 are positive, and we set

$$a = \min \left\{ \sqrt{\hat{\psi}''(0)}, \sqrt{\hat{\psi}''(1)} \right\}.$$

Depending on the parameters $\lambda$ and $\mu$, we partition $Q$ into the inner neighborhood $Q^{(\mu\lambda)}_{\text{int}}$ of $\Gamma$, into the matching region $Q^{(\mu\lambda)}_{\text{match}}$ and into the outer region $Q^{(\mu\lambda)}_{\text{out}}$. These sets are
defined by

\[ Q_{\text{inn}}^{(\mu\lambda)} = \left\{ (t, \eta, \xi) \in \mathcal{U}_\delta \mid |\xi| < \frac{3}{2} \frac{\ln \mu}{a} \right\}, \]

\[ Q_{\text{match}}^{(\mu\lambda)} = \left\{ (t, \eta, \xi) \in \mathcal{U}_\delta \mid \frac{3}{2} \frac{\ln \mu}{a} \leq |\xi| \leq \frac{3}{2} \frac{\ln \mu}{a} \right\}, \quad (2.34) \]

\[ Q_{\text{out}}^{(\mu\lambda)} = Q \setminus (Q_{\text{inn}}^{(\mu\lambda)} \cup Q_{\text{match}}^{(\mu\lambda)}). \]

We always assume that the parameters \(\lambda\) and \(\mu\) satisfy \(0 < \lambda \leq \lambda_0\) and \(0 < \mu \leq \mu_0\), where \(\lambda_0, \mu_0\) are fixed constants satisfying

\[ \mu_0 \leq e^{-2}, \quad \frac{3(\mu_0 \lambda_0)^{1/2} \ln \mu_0}{a} < \delta. \]

The first condition is imposed for purely technical reasons and guarantees that the function \(\mu \mapsto \frac{\ln \mu}{\mu} a\) is increasing, the second condition guarantees that \(Q_{\text{inn}}^{(\mu\lambda)}, Q_{\text{match}}^{(\mu\lambda)} \subset \mathcal{U}_\delta\) and that \(Q_{\text{out}}^{(\mu\lambda)} \cap \mathcal{U}_\delta\) is a nonempty, relatively open subset of \(\mathcal{U}_\delta\).

By (2.18), the function \(\hat{u} : \mathcal{U} \to \mathbb{R}^3\) is continuous at every point \((t, \eta) \in \Gamma\), but the first and higher derivatives of \(\hat{u}\) in the direction of the normal vector \(n(t, x)\) can jump across \(\Gamma\). For these jumps we write

\[ u^*(t, \eta) = [\partial_n \hat{u}(t, \eta)], \quad (2.35) \]

\[ a^*(t, \eta) = [\partial_n^2 \hat{u}(t, \eta)]. \quad (2.36) \]

We also set

\[ c_1 = \int_0^1 \sqrt{2 \psi(\vartheta)} \, d\vartheta. \quad (2.37) \]

**Theorem 2.3** *Suppose that the double well potential \(\hat{\psi} \in C^5(\mathbb{R})\) satisfies*

\[ \hat{\psi}(r) > 0, \quad \text{for } 0 < r < 1, \]

\[ \hat{\psi}(r) = \hat{\psi}'(r) = 0, \quad \text{for } r = 0, 1, \]

\[ a = \min \left\{ \sqrt{\hat{\psi}''(0)}, \sqrt{\hat{\psi}''(1)} \right\} > 0. \quad (2.38) \]

*Moreover, suppose that \(\hat{\psi}\) satisfies the symmetry condition*

\[ \hat{\psi}\left(\frac{1}{2} - \zeta\right) = \hat{\psi}\left(\frac{1}{2} + \zeta\right), \quad \zeta \in \mathbb{R}. \quad (2.39) \]

*Assume that there is a solution \(\Gamma\) of the evolution problem (2.14), (2.15) with \(s_0 = s_0(\hat{T}, \kappa_G, \lambda^{1/2}) : \Gamma \to \mathbb{R}\) given by*

\[ s_0 = \frac{c}{c_1} \left( -\varpi : (\hat{T}) + \lambda^{1/2} c_1 \kappa_G \right), \quad (2.40) \]

*and with \(s_1 = s_1(\hat{u}, \hat{T}, \hat{T}, S_0, S_1, \lambda^{1/2}) : \Gamma \to \mathbb{R}\) defined by*

\[ s_1 = s_{10} + \lambda^{1/2} s_{11} = s_{10}(\hat{T}, \hat{T}, S_0, S_1) + \lambda^{1/2} s_{11}(\hat{u}, S_0), \quad (2.41) \]

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\[ s_{10} = \frac{c}{c_1} \left( \frac{\pi}{\psi''(S)} \right) - \int_{-\infty}^{\infty} S_1 S_0' \, d\zeta \right) + \frac{1}{c_1} \pi \left( \frac{\pi}{\psi''(S)} \right) \int_{-\infty}^{\infty} S_1 S_0' \, d\zeta + \frac{1}{2} \int_{-\infty}^{\infty} \pi''(S) S_1 S_0' \, d\zeta, \quad (2.42) \]

\[ s_{11} = -\frac{c}{c_1} \pi \left( D e(a^* \otimes n + \nabla T u^*) \right) \int_{-\infty}^{\infty} S_0(\zeta) S_0(-\zeta) \, d\zeta. \quad (2.43) \]

In (2.44) and (2.45), we have \( S_0 = S_0(\zeta) \) and \( S_1 = S_1(t, \eta, \zeta) \), for all other functions the argument is \((t, \eta)\). The positive constant \( c \) is defined in (1.3). The notations \([\cdot]\) and \(\langle \cdot \rangle\) are introduced in Section 2.1. In particular, we have

\[ \langle \frac{\pi}{\psi''(S)} \rangle = \frac{1}{2} \left( \frac{\pi}{\psi''(1)} + \frac{\pi}{\psi''(0)} \right). \]

With these functions the normal speed \( s(t, \eta) \) of \( \Gamma(t) \) at \( \eta \in \Gamma(t) \) is thus given by

\[ s(t, \eta) = s_0(t, \eta) + \mu^2 s_1(t, \eta, \lambda^2) = s_0(t, \eta) + \mu^{1/2} \left( s_{10}(t, \eta) + \lambda^{1/2} s_{11}(t, \eta) \right). \quad (2.44) \]

We assume moreover that the solution \( \Gamma \) is a \( C^3 \)–manifold and that the functions \( \bar{u} \) and \( \hat{u} \) defined by the evolution problem satisfy \( \bar{u} \in C^4(\gamma \cup \gamma', \mathbb{R}^3) \), \( \hat{u} \in C^3(\gamma \cup \gamma', \mathbb{R}^3) \) and that \( \bar{u} \) has \( C^4 \)–extensions, \( \hat{u} \) has \( C^3 \)–extensions from \( \gamma \) to \( \gamma \cup \Gamma \) and from \( \gamma' \) to \( \gamma' \cup \Gamma \). For the given right hand side of (1.1) we assume that \( b \in C^1(\bar{Q}) \).

Under these assumptions there is an approximate solution \((u(\mu), T(\mu), S(\mu))\) of the Allen-Cahn model \( (1.1) \) – \( (1.5) \), for which \( \Gamma \) is the level set

\[ \Gamma = \left\{ (t, x) \in Q \mid S(\mu)(t, x) = \frac{1}{2} \right\}, \quad (2.45) \]

and which satisfies the equations

\[ -\text{div} x T(\mu) = b + f_1(\mu, \lambda), \quad (2.46) \]

\[ T(\mu) = D \left( e(\nabla x u(\mu)) - \pi S(\mu) \right), \quad (2.47) \]

\[ \partial_t S(\mu) + \frac{c}{(\mu, \lambda)\overline{2}} \left( \partial_\Sigma W(\nabla x u(\mu), S(\mu)) + \frac{1}{\mu^{1/2}} \psi'(S(\mu)) - \mu^{1/2} \lambda S(\mu) \right) = f_2(\mu, \lambda), \quad (2.48) \]

\[ u(\mu)(t, x) = U(t, x), \quad (t, x) \in [t_1, t_2] \times \partial \Omega, \quad (2.49) \]

\[ \partial_{n, \text{out}} S(\mu)(t, x) = f_3(\mu, \lambda), \quad (t, x) \in [t_1, t_2] \times \partial \Omega, \quad (2.50) \]

where to the right hand sides \( f_1(\mu, \lambda), \ldots, f_3(\mu, \lambda) \) there exist nonnegative constants \( K_1, \ldots, K_5 \) such that for all \( \mu \in (0, \mu_0) \) and all \( \lambda \in (0, \lambda_0) \]

\[ \| f_1(\mu, \lambda) \|_{L^\infty(Q_{\text{inn}}^{(\mu, \lambda)}, \Omega_{\text{match}}^{(\mu, \lambda)})} \leq \| \ln \mu \|^{2} \left( \frac{\mu}{\lambda} \right)^{1/2} K_1, \quad (2.51) \]

\[ \| f_1(\mu, \lambda) \|_{L^\infty(Q_{\text{out}}^{(\mu, \lambda)})} \leq \mu^{1/2} K_2, \quad (2.52) \]

\[ \| f_2(\mu, \lambda) \|_{L^\infty(Q_{\text{inn}}^{(\mu, \lambda)}, \Omega_{\text{match}}^{(\mu, \lambda)})} \leq \| \ln \mu \|^{2} \left( \frac{\mu}{\lambda} \right)^{1/2} K_3, \quad (2.53) \]

\[ \| f_2(\mu, \lambda) \|_{L^\infty(Q_{\text{out}}^{(\mu, \lambda)})} \leq \frac{\mu}{\lambda^{1/2}} K_4, \quad (2.54) \]

\[ \| f_3(\mu, \lambda) \|_{L^\infty(\partial \Omega)} \leq \mu^{1/2} K_5, \quad (2.55) \]
In the neighborhood $Q_{\text{inn}}^{(\mu\lambda)}$ of $\Gamma$ the order parameter in the approximate solution is of the form

$$S^{(\mu)}(t,x) = S_0\left(\xi_{(\mu\lambda)^{1/2}}\right) + \mu^{1/2}S_1\left(t,\eta,\xi_{(\mu\lambda)^{1/2}}\right) + \mu S_2\left(t,\eta,\xi_{(\mu\lambda)^{1/2}}\right),$$  \hspace{1cm} (2.56)

where the monotonically increasing transition profile $S_0 : \mathbb{R} \rightarrow \mathbb{R}$ and the function $S_1 : \Gamma \times \mathbb{R} \rightarrow \mathbb{R}$ are given as solution of the coupled problem (2.26) – (2.32), and where $S_2 : \Gamma \times \mathbb{R} \rightarrow \mathbb{R}$ satisfies $S_2(t,\eta,0) = 0$ and

$$|S_2(t,\eta,\zeta)| \leq C(1 + |\zeta|), \hspace{1cm} \text{for} \hspace{0.2cm} (t,\eta,\zeta) \in \Gamma \times \mathbb{R},$$  \hspace{1cm} (2.57)

with a constant $C$ independent of $(t,\eta,\zeta)$.

We mention that the positive constant $c$ in (1.3) does not play a major role in the analysis and could be replaced by 1. We refrain from replacing it to show how $c$ appears in the kinetic relation.

The proof of this theorem forms the content of Sections 3 – 5. We remark that the symmetry assumption (2.39) for the double well potential $\hat{\psi}$ serves to simplify the computations in the derivation of the asymptotic solution. Without this assumption the term $s_1$ in the kinetic relation (2.44) would contain other terms in addition to the terms $s_{10}$ and $s_{11}$ given in (2.42) and (2.43).

The regularity properties of $\Gamma$ and of $\hat{u}$, $\hat{u}$ are of course not independent, since $\hat{u}$ and $\hat{u}$ are solutions of the elliptic transmission problems (2.16) – (2.20) and (2.21) – (2.25), respectively. Therefore the regularity theory of elliptic equations shows that $\hat{u}$ and $\hat{u}$ automatically have the differentiability properties assumed in the theorem if the manifold $\Gamma$ and the right hand side $b$ are sufficiently smooth.

Since by definition of $Q_{\text{inn}}^{(\mu\lambda)}$ and $Q_{\text{match}}^{(\mu\lambda)}$ in (2.34) we have

$$\text{meas}(Q_{\text{inn}}^{(\mu\lambda)} \cup Q_{\text{match}}^{(\mu\lambda)}) \leq C_3(\mu\lambda)^{1/2} |\ln \mu|,$$

we immediately obtain from (2.51) – (2.54) the following

**Corollary 2.4** There are constants $K_6$, $K_7$ such that for all $0 < \mu \leq \mu_0$ and all $0 < \lambda \leq \lambda_0$

$$\|f_1^{(\mu\lambda)}\|_{L^1(Q)} \leq |\ln \mu|^3 \mu K_6,$$  \hspace{1cm} (2.58)

$$\|f_2^{(\mu\lambda)}\|_{L^1(Q)} \leq \frac{|\ln \mu|^3 \mu}{\lambda^{1/2}} K_7.$$  \hspace{1cm} (2.59)

The leading term $s_0$ given in (2.40) can be written in a more common and more general form. To give this form, we need a result on the jump of the Eshelby tensor. The Eshelby tensor to the solution $(\hat{u}, \hat{T})$ of the transmission problem (2.16) – (2.20) is defined by

$$\hat{C}(\nabla_x \hat{u}, \hat{S}) = \psi_{\mu}(\varepsilon(\nabla_x \hat{u}), \hat{S})I - (I + \nabla_x \hat{u})^T \hat{T},$$  \hspace{1cm} (2.60)

where $I \in \mathbb{R}^3$ is the unit matrix and where

$$\psi_{\mu}(\varepsilon, S) = W(\varepsilon, S) + \frac{1}{\mu^{1/2}} \hat{\psi}(S)$$  \hspace{1cm} (2.61)

is that part of the free energy $\psi_{\mu}^*$ defined in (1.9) without gradient term. The last term on the right hand side of (2.60) is a matrix product. We use the standard convention to denote the matrix product of two matrices $A \in \mathbb{R}^{k \times m}$ and $B \in \mathbb{R}^{m \times \ell}$ by $AB \in \mathbb{R}^{k \times \ell}$.
Lemma 2.5  Let \((\hat{u}, \hat{T})\) be the solution of the transmission problem (2.16) – (2.20) and let \(n\) be a unit normal vector field to \(\Gamma(t)\). Then the jump \([\hat{C}]\) of the Eshelby tensor to \((\hat{u}, \hat{T})\) across \(\Gamma\) satisfies
\[
n \cdot [\hat{C}] n = \frac{1}{\mu^{1/2}} [\hat{\psi}(\hat{S})] - \varpi : \langle \hat{T} \rangle.
\] (2.62)

This result is known [2]. In [7] it is stated as equation (3.4) and proved on pages 154, 155.

Corollary 2.6  The leading term \(s_0\) of the kinetic relation defined in (2.41) satisfies
\[
s_0 = \frac{c}{c_1} \left( n \cdot [\hat{C}] n + \lambda^{1/2} c_1 \kappa_{\Gamma} \right).
\] (2.63)

This corollary follows immediately from (2.62), since by assumption (2.38) we have 
\([\hat{\psi}(\hat{S})] = \hat{\psi}(1) - \hat{\psi}(0) = 0\), which implies that 
\(n \cdot [\hat{C}] n = -\varpi : \langle \hat{T} \rangle\).

2.4  Consequences for numerical simulations

In this section we discuss the consequences of Theorem 2.3 for numerical simulations of interfaces with small interface energy.

In many functional materials the phase interfaces consist only of a few atomic layers. For interfaces with such small width mathematical models with sharp interface are appropriate. We therefore base the following considerations on the hypothesis that the propagation speed of the interface in the sharp interface model is a good approximation to the propagation speed of the interface in the real material. The model error of the Allen-Cahn model is then the difference of the propagation speed of the sharp interface and the propagation speed of the diffuse interface in the phase field model. The parameters \(\mu\) and \(\lambda\) in the Allen-Cahn model should be chosen such that this model error is small and such that numerical simulations based on the Allen-Cahn model are effective.

To make this precise we must first determine the sharp interface model to be used. The model consists of the transmission problem (2.16) – (2.20) combined with a kinetic relation. To find this relation, one proceeds in the usual way and uses that by the second law of thermodynamics the Clausius-Duhem inequality
\[
\partial_t \psi_{\text{sharp}} + \text{div}_x q_{\text{sharp}} \leq \hat{u}_t \cdot b
\]
must be satisfied to impose restrictions on the form of the kinetic relation. Here \(\psi_{\text{sharp}}\) denotes the free energy in the sharp interface problem and \(q_{\text{sharp}}\) is the flux of the free energy. We use the standard free energy and flux
\[
\psi_{\text{sharp}}(\varepsilon(\nabla_x \hat{u}), \hat{S}) = W(\varepsilon(\nabla_x \hat{u}), \hat{S}) + \lambda^{1/2} c_1 \int_{\Gamma(t)} d\sigma,
\] (2.64)
\[
q_{\text{sharp}}(\hat{T}, \hat{S}) = -\hat{T} \cdot \hat{u}_t.
\]

The last term on the right hand side of (2.64) is the interface energy, hence \(\lambda^{1/2} c_1\) is the interface energy density. It is well known that if \((\hat{u}(t), \hat{T}(t))\) is a solution of the transmission problem (2.16) – (2.20) at time \(t\) and if the interface \(\Gamma(t)\) in this problem
moves with the given normal speed \( s_{\text{sharp}}(t, x) \) at \( x \in \Gamma(t) \), then the Clausius-Duhem inequality holds if and only if the inequality

\[
s_{\text{sharp}}(t, x) \left( n(t, x) \cdot [\hat{C}] (t, x) n(t, x) + \lambda^{1/2} c_1 \kappa (t, x) \right) \geq 0
\]  

(2.65)

is satisfied at every point \( x \in \Gamma(t) \). A proof of this well known result is given in [3], however only for the case where \( \lambda = 0 \) in (2.64). The proof can be readily generalized to the case \( \lambda > 0 \).

A simple linear kinetic relation, for which (2.65) obviously holds, is

\[
s_{\text{sharp}} = \frac{c}{c_1} \left( n \cdot [\hat{C}] n + \lambda^{1/2} c_1 \kappa \right).
\]

(2.66)

The sharp interface problem thus consists of the transmission problem (2.16) – (2.20) combined with the kinetic relation (2.66). For this problem the Clausius-Duhem inequality is satisfied.

We can now define the model error. To this end let \( (u^{(\mu)}, T^{(\mu)}, S^{(\mu)}) \) be the asymptotic solution in the domain \( Q = [t_1, t_2] \times \Omega \) constructed in Theorem 2.3, where by (2.45), the manifold \( \Gamma \) is the level set \( \{S^{(\mu)} = \frac{1}{2}\} \). Let \( t \in [t_1, t_2] \) be a fixed time and let \( (u^{(AC)}, T^{(AC)}, S^{(AC)}) \) be the exact solution of the Allen-Cahn model (1.1) – (1.3) in the domain \( [\hat{t}, t_2] \times \Omega \), which satisfies the boundary and initial conditions

\[
\begin{align*}
    u^{(AC)}(t, x) &= U(t, x), \quad (t, x) \in [\hat{t}, t_2] \times \partial \Omega, \\
    \partial_n S^{(AC)}(t, x) &= f_3^{(\mu, \lambda)}(t, x), \quad (t, x) \in [\hat{t}, t_2] \times \partial \Omega, \\
    S^{(AC)}(\hat{t}, x) &= S^{(\mu)}(\hat{t}, x), \quad x \in \Omega,
\end{align*}
\]

(2.67) – (2.69)

where \( f_3^{(\mu, \lambda)} \) is the right hand side of (2.50). The level set of the order parameter \( S^{(\mu)}_{AC} \) is denoted by

\[
\Gamma_{AC} = \left\{ (t, x) \in Q : S^{(\mu)}_{AC}(t, x) = \frac{1}{2} \right\},
\]

(2.70)

Let \( \Gamma_{\text{sharp}} \subseteq Q \) be the sharp interface in the solution of the sharp interface problem (2.16) – (2.20), (2.66), which satisfies the initial condition

\[
\Gamma_{\text{sharp}}(\hat{t}) = \Gamma(\hat{t}).
\]

(2.71)

The normal speeds of the different surfaces are

\[
s = s^{(\mu, \lambda)} = N(\Gamma^{(\mu, \lambda)}), \quad s_{AC} = s^{(\mu, \lambda)}_{AC} = N(\Gamma^{(\mu, \lambda)}_{AC}), \quad s_{\text{sharp}} = N(\Gamma_{\text{sharp}}),
\]

where \( N \) is the normal speed operator introduced at the beginning of Section 2.2. Of course, \( s_{\text{sharp}} \) is given by (2.66). Note that the functions \( s^{(\mu, \lambda)}(\hat{t}), s^{(\mu, \lambda)}_{AC}(\hat{t}), s_{\text{sharp}}(\hat{t}) \) are defined on the same set, since the initial condition (2.69) and (2.71) together imply \( \Gamma_{AC}(\hat{t}) = \Gamma(\hat{t}) = \Gamma_{\text{sharp}}(\hat{t}) \).

**Definition 2.7** We call the function \( \mathcal{E} = \mathcal{E}^{(\mu, \lambda)}(\hat{t}) : \Gamma(\hat{t}) \to \mathbb{R} \) defined by

\[
\mathcal{E} = s_{AC}(\hat{t}) - s_{\text{sharp}}(\hat{t})
\]

(2.72)

the model error of the Allen-Cahn model at time \( \hat{t} \) to the parameters \( \mu \) and \( \lambda \).
We can now discuss the choice of the parameters $\mu$ and $\lambda$. Since (2.66) coincides with the leading term $s_0$ in the asymptotic expansion (2.44) of the kinetic relation of the Allen-Cahn model, which is seen from (2.63), we have

$$s_{\text{sharp}} = s_0.$$  

(2.73)

Therefore (2.44) yields

$$E = s_{AC} - s_{\text{sharp}} = s_{AC} - s_0 = (\mu s_{AC} - s) + (s - s_0) = (\mu s_{AC} - s) + \mu^{1/2}(s_{10} + \lambda^{1/2}s_{11}).$$  

(2.74)

The difference $s_{AC} - s$ between the propagation speeds of the exact solution and the asymptotic solution tends to zero for $\mu \to 0$ faster than the term $\mu^{1/2}s_{10}$, and the convergence is uniform with respect to $\lambda$. This is the basic result, which allows to discuss the optimal choice of $\mu$ and $\lambda$. The precise result is

**Theorem 2.8** There is a constant $C_\epsilon > 0$ such that for all $0 < \mu \leq \mu_0$ and all $0 < \lambda \leq \lambda_0$ we have the estimate

$$\|s_{AC}(\hat{t}) - s(\hat{t})\|_{L^2(\Gamma(\hat{t}))} \leq C_\epsilon |\ln \mu|^3 \mu.$$  

(2.75)

The proof of this theorem is given in Section 6. (2.74) and (2.75) together yield

$$\|\mathcal{E}(\mu\lambda)\|_{L^2(\Gamma(\hat{t}))} \leq C \mu^{1/2},$$  

(2.76)

with a constant $C$, which can be chosen independently of $\lambda$. By this inequality, $\mu^{1/2}$ controls the model error. Therefore we write $F = \mu^{1/2}$ and call $F$ the error parameter. Moreover, since $\lambda^{1/2} c_1$ is the interface energy density, we call $E = \lambda^{1/2}$ the interface energy parameter. Also, since by (2.56) the interface width is proportional to $(\mu\lambda)^{1/2}$, we call $B = (\mu\lambda)^{1/2}$ the interface width parameter. These three parameters and the propagation speed $s_{AC}$ are connected by the fundamental relations

$$B = EF,$$  

(2.77)

$$s_{AC} = \frac{c}{c_1} n \cdot [\hat{C}] n + c\kappa E + \mathcal{E}[E, F],$$  

(2.78)

$$\|\mathcal{E}[E, F]\|_{L^2(\Gamma(\hat{t}))} \leq CF,$$  

(2.79)

where we use the notation $\mathcal{E}[E, F] = \mathcal{E}(\mu\lambda)$. The first equation is an immediate consequence of the definition of the parameters, the second is obtained by insertion of (2.66) into (2.72), and the last inequality is just a restatement of (2.76).

Now assume that we want to use a phase field model to numerically simulate the propagation of a phase interface. In such a simulation the numerical effort is proportional to $h^{-p}$, where $h$ denotes the grid spacing and where the power $p > 1$ depends on whether we want to simulate a problem in 2–d or in 3–d and it depends on the numerical scheme we use. In order for the simulation to be precise, we must guarantee that the model error and the numerical error are small. To make the numerical error small, we must choose the grid spacing $h$ small enough to resolve the transition of the order parameter across the interface, which means that we must choose $h < B$, hence we have $h^{-p} > B^{-p}$. 

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Therefore we see that the numerical effort of a simulation based on a phase field model is measured by the number $B^{-p}$. We call the number

$$e_{\text{num}} = B^{-p}$$

the parameter of numerical effort. For a simulation based on the Allen-Cahn model we see from (2.77) that the numerical effort is

$$e_{\text{num}} = (EF)^{-p}.$$  

Assume that the interface, which we want to simulate with the Allen-Cahn model, has very small interface energy density. Such interfaces are common in metallic or functional materials. For such materials the interface energy parameter $E$ is small. To make the model error small, we must also choose the error parameter $F$ small, which means that the numerical effort parameter $e_{\text{num}} = (EF)^{-p}$ is very large as a product of two large numbers $E^{-p}$ and $F^{-p}$.

To be more specific, we consider an interface without interface energy, which means that the free energy $\psi_{\text{sharp}}$ does not contain the last term on the right hand side of (2.64). From (2.66) we see that the propagation speed of the sharp interface with zero interface energy density is

$$s_{\text{sharp}} = c_{c1} \cdot [\hat{C}].$$

From this equation and from (2.78) we see that in this case the total model error, which we denote by $E_{\text{total}}$, is

$$E_{\text{total}} = s_{AC} - s_{\text{sharp}} = c_{\text{sharp}} E + E[F].$$

This means that the term $c_{\text{sharp}} E$ is now part of the total model error. This term does not vanish identically, since we cannot set $\lambda = 0$ in the Allen-Cahn equation (1.3). Instead the values of $\lambda$ and of $E = \lambda^{1/2}$ must be positive.

If we prescribe the $L^2$-norm $E_{L^2} = \|E_{\text{total}}\|_{L^2(\Gamma(t))}$ of the total model error, we must therefore choose the parameters $E$ and $F$ such that

$$c_{\text{sharp}} \|E_{\text{total}}\|_{L^2(\Gamma(t))} E + \|E[F]\|_{L^2(\Gamma(t))} \leq E_{L^2},$$

where the second condition is imposed by the requirement to make the numerical effort $e_{\text{num}} = (EF)^{-p}$ as small as possible. To discuss this optimization problem, we assume first that the term $s_{10}$ in the asymptotic expansion (2.44) of the kinetic relation of the Allen-Cahn model is not identically equal to zero. In this case we conclude from (2.74) and (2.75) by the inverse triangle inequality that for sufficiently small $\lambda^{1/2} = E$ and for sufficiently small $\mu^{1/2} = F$

$$\|E[F]\|_{L^2(\Gamma(t))} \geq \mu^{1/2} s_{10} + \frac{(\mu \lambda)^{1/2}}{2} s_{11} + \frac{(s_{\text{AC}} - s)^{1/2}}{2} L^2(\Gamma(t)) - \|s_{\text{AC}} - s\|_{L^2(\Gamma(t))} \geq \mu^{1/2} s_{10} + \frac{(\mu \lambda)^{1/2}}{2} s_{11} + \frac{(s_{\text{AC}} - s)^{1/2}}{2} L^2(\Gamma(t)) - C_\varepsilon \|s_{11}\|_{L^2(\Gamma(t))} - C_\varepsilon \|s_{\text{AC}} - s\|_{L^2(\Gamma(t))}$$

where the second condition is imposed by the requirement to make the numerical effort $e_{\text{num}} = (EF)^{-p}$ as small as possible.
This inequality and (2.80) imply that the solution \((E, F)\) of the optimization problem (2.80), (2.81) satisfies
\[
F \leq \frac{2}{\|s_{10}\|_{L^2(\Gamma(\hat{t}))}} \| \mathcal{E}[E, F]\|_{L^2(\Gamma(\hat{t}))} \leq \frac{2}{\|s_{10}\|_{L^2(\Gamma(\hat{t}))}} \mathcal{E}_{L^2} \quad \text{and} \quad E \leq \frac{1}{c\|\kappa_T\|_{L^2(\Gamma(\hat{t}))}} \mathcal{E}_{L^2}.
\]
From this result we obtain

**Corollary 2.9** Let \(\mathcal{E}_{\text{max}}\) denote the total model error of the Allen-Cahn model in the simulation of an interface without interface energy. If the term \(s_{10}\) in the asymptotic expansion (2.44) of the kinetic relation of the Allen-Cahn model is not identically equal to zero, then the interface width \(B\) satisfies
\[
B = EF \leq \frac{2}{c\|s_{10}\|_{L^2(\Gamma(\hat{t}))}\|\kappa_T\|_{L^2(\Gamma(\hat{t}))}} \mathcal{E}_{L^2}^2.
\]
In a numerical simulation of an interface without interface energy the parameter of numerical effort satisfies
\[
\epsilon_{\text{num}} \geq \left( \frac{c\|s_{10}\|_{L^2(\Gamma(\hat{t}))}\|\kappa_T\|_{L^2(\Gamma(\hat{t}))}}{2E_{L^2}^2} \right)^p
\]
with a power \(p > 1\) depending on the space dimension and the numerical method used.

The interface width thus decreases with the square of the model error. Since the time step in a simulation must be decreased when the grid spacing \(h\) in \(x\)-direction is decreased, the number \(p\) can be larger than 4 in a three dimensional simulation. From (2.83) we thus see that the numerical effort grows very rapidly when the required accuracy is increased. The Allen-Cahn model is therefore ineffective when used to accurately simulate interfaces with low interface energy.

If the term \(s_{10}\) vanishes identically, then the same considerations show that instead of (2.82) and (2.83) we would have \(B = O(\mathcal{E}_{L^2}^{1/2})\) and \(\epsilon_{\text{num}} \geq C\mathcal{E}_{L^2}^{1/2p}\). The numerical effort would still grow fast when the required accuracy is increased, though less fast than for \(s_{10} \neq 0\). However, a close investigation of the terms in the definition (2.42) of \(s_{10}\), which we do not present here, shows that only in very exceptional situations one can expect that \(s_{10}\) vanishes identically.

In Corollary 2.9 we assumed that the mesh is globally refined. Of course, one can improve the effectiveness of simulations by using local mesh refinement in the neighborhood of the interface. We do not discuss this question of numerical analysis here, but Corollary 2.9 in fact shows that adaptive mesh refinement and other advanced numerical techniques are needed to make precise simulations of interfaces with small energy based on the Allen-Cahn model effective.

**Comparison to the hybrid phase field model** With Corollary 2.9 we can refine the comparison given in [8] of the Allen-Cahn model and another phase field model, which we call the hybrid model. The hybrid model was introduced and discussed in [3, 4, 5, 6, 8]. By formal construction of asymptotic solutions we showed in [8] the following result:

Let \(B_{AC}(\mathcal{E}_{L^2})\) and \(B_{hyb}(\mathcal{E}_{L^2})\) be the interface widths in the Allen-Cahn model and the
hybrid model, respectively, which result when the model parameters are adjusted to
model an interface without interface energy with the total model error \( \mathcal{E}_{L^2} \). Then we
have for \( \mathcal{E}_{L^2} \to 0 \) that

\[
B_{hyb}(\mathcal{E}_{L^2}) = O(\mathcal{E}_{L^2}), \quad B_{AC}(\mathcal{E}_{L^2}) = o(1)O(\mathcal{E}_{L^2}) = o(\mathcal{E}_{L^2}).
\]

The Landau symbol \( o(1) \) denotes terms, which tend to zero for \( \mathcal{E}_{\text{max}} \to 0 \). This result
was obtained under the assumption that estimates corresponding to (2.75) hold true for
both models, without proving this assumption.

To achieve a prescribed small value \( \mathcal{E}_{L^2} \) of the total model error we must therefore
choose the interface width in the Allen-Cahn model smaller than in the hybrid model.
Consequently, the hybrid model is numerically more effective, but how much more de-
pends on the rate of decay of the \( o(\mathcal{E}_{L^2}) \) term in the result for the Allen-Cahn model. In
[8] we could not determine this decay rate, since the asymptotic solution constructed in
[8] for the Allen-Cahn model was only of first order.

Corollary 2.9 yields this decay rate. From the result for the hybrid model and from
Corollary 2.9 we thus obtain for the parameters \( e_{\text{num}}^{hyb} \) and \( e_{\text{num}}^{AC} \) of the hybrid model and
the Allen-Cahn model, respectively, that

\[
e_{\text{num}}^{hyb} \leq C\mathcal{E}_{L^2}^{-p}, \quad e_{\text{num}}^{AC} \geq C\mathcal{E}_{L^2}^{-2p},
\]

which shows that when the prescribed error \( \mathcal{E}_{L^2} \) is small, the hybrid model can be quite
considerably more effective in numerical simulations of interfaces with low interface energy
or no interface energy than the Allen-Cahn model.

2.5 The jump of solutions of the transmission problems

In this section we prove some results on the jumps of the solutions \((\hat{u}, \hat{T})\) and \((\tilde{u}, \tilde{T})\)
of the transmission problems (2.16) – (2.20) and (2.21) – (2.25), which we need in the
following sections.

We define a scalar product \( \alpha : \beta \) on \( S^3 \) by \( \alpha : \beta = \alpha : (D\beta) \), for \( \alpha, \beta \in S^3 \). For a
unit vector \( n \in \mathbb{R}^3 \) let a linear subspace of \( S^3 \) be given by

\[
S^3_n = \left\{ \frac{1}{2}(\omega \otimes n + n \otimes \omega) \mid \omega \in \mathbb{R}^3 \right\},
\]

let \( P_n : S^3 \to S^3 \) be the projector onto \( S^3_n \), which is orthogonal with respect to the scalar
product \( \alpha : \beta \) and let \( Q_n = I - P_n \).

**Lemma 2.10** Let \( \omega^* \in \mathbb{R}^3 \) be a vector. This vector satisfies \( D(\varepsilon(\omega^* \otimes n) - \overline{\gamma}) \) \( n = 0 \) if
and only if \( \varepsilon(\omega^* \otimes n) = P_n\overline{\gamma} \) holds.

This lemma is proved in [7, Lemma 2.2].

**Lemma 2.11** Let \((\hat{u}, \hat{T})\) be a solution of the transmission problem (2.16) – (2.20). Assume
that \( \hat{u} \) is continuous in \( Q \) and that the limits \( (\nabla_x \hat{u})^{(\pm)} \) exist and define continuous
extensions of \( \nabla_x \hat{u} \) from the set \( \gamma' \) to \( \gamma' \cup \Gamma \) and from the set \( \gamma \) to \( \gamma \cup \Gamma \), respectively. Then we have

\[
[\varepsilon(\nabla_x \hat{u})] = \varepsilon(u^* \otimes n), \quad [\hat{T}] = D(\varepsilon(u^* \otimes n) - \overline{\gamma}), \quad (2.85)
\]

\[
[\varepsilon(\nabla_x \tilde{u})] = P_n\overline{\gamma}, \quad [\tilde{T}] = -DQ_n\overline{\gamma}, \quad (2.86)
\]
Thus, by (2.17),
\[
D(\varepsilon(u^\ast \otimes n) - \bar{\gamma}) = D(\varepsilon(\nabla_x \hat{u}) - \bar{\gamma}[S]) = [D(\varepsilon(\nabla_x \hat{u}) - \bar{\gamma}S)] = [\bar{T}].
\]
This proves (2.85). From (2.19) and (2.85) we infer that
\[
0 = [\bar{T}]n = \left(D(\varepsilon(u^\ast \otimes n) - \bar{\gamma})\right)n,
\]
so that \(\varepsilon(\nabla_x \hat{u}) = \varepsilon(u^\ast \otimes n) = P_n\bar{\gamma},\) by Lemma 2.10. Therefore we find
\[
[\bar{T}] = D(\varepsilon(\nabla_x \hat{u}) - \bar{\gamma}) = D(P_n\bar{\gamma} - \bar{\gamma}) = -DQ_n\bar{\gamma},
\]
which proves (2.86).

\textbf{Lemma 2.12} Let \((\hat{u}, \bar{T})\) and \((\hat{u}, \bar{T})\) be solutions of the transmission problems (2.16) – (2.20) and (2.21) – (2.24), respectively. Assume that \(\hat{u}\) and \(\hat{u}\) are continuous in \(Q\) and that the limits \(\nabla_x \hat{u})^{(\pm)}\) and \(\nabla_x \hat{u})^{(\pm)}\) exist and define continuous extensions of \(\nabla_x \hat{u}\) and of \(\nabla_x \hat{u}\) from the set \(\gamma'\) to \(\gamma' \cup \Gamma\) and from the set \(\gamma\) to \(\gamma \cup \Gamma\), respectively. Then we have
\[
\partial_n \hat{u} = \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast, \quad \nabla_x \hat{u} = \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast \otimes n. \tag{2.88}
\]
\textbf{Proof}: The decomposition (2.8) yields
\[
\nabla_x \hat{u} = (\partial_n \hat{u} \otimes n + \nabla_x \gamma) = (\partial_n \hat{u} \otimes n + \nabla_x \gamma) = \partial_n \hat{u} \otimes n. \tag{2.89}
\]
From this equation and from (2.22), (2.24) we infer
\[
0 = [\bar{T}]n = \left(D(\varepsilon(\nabla_x \hat{u}) - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]\bar{\gamma})\right)n = \left(D(\varepsilon(\partial_n \hat{u}) \otimes n) - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]\bar{\gamma}\right)n.
\]
Thus, Lemma 2.10, the linearity of the projector \(P_n\) and (2.85), (2.86) imply
\[
\varepsilon(\partial_n \hat{u}) \otimes n = \left[\frac{\bar{\gamma}}{\psi''(S)}\right]P_n\bar{\gamma} = \left[\frac{\bar{\gamma}}{\psi''(S)}\right]\varepsilon(u^\ast \otimes n) = \varepsilon\left(\left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast \otimes n\right),
\]
whence
\[
\left(\partial_n \hat{u} - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast\right) \otimes n + n \otimes \left(\partial_n \hat{u} - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast\right) = 0.
\]
We multiply this equation from the right with \(n\) and obtain
\[
\left(\partial_n \hat{u} - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast\right) + n\left(\partial_n \hat{u} - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast\right) \cdot n = 0,
\]
which means that \(\partial_n \hat{u} - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast\) is a multiple of \(n\). Scalar multiplication of the last equation with \(n\) yields \((\partial_n \hat{u} - \left[\frac{\bar{\gamma}}{\psi''(S)}\right]u^\ast) \cdot n = 0,\) from which we now conclude that the first equation in (2.88) holds. The second equation is obtained from (2.89).
3 The asymptotic solution

This section forms the first part of the proof Theorem 2.3. We state in this section the form of the asymptotic solution \( (u^{(\mu)}, T^{(\mu)}, S^{(\mu)}) \). The proof continues in Section 4 where we study properties of the functions \( S_0, \ldots, S_2 \) appearing in the asymptotic solution. In Section 5 we use these properties to show that \( (u^{(\mu)}, T^{(\mu)}, S^{(\mu)}) \) is an asymptotic solution by verifying that the estimates (2.51) – (2.54) hold. This concludes the proof of Theorem 2.3.

3.1 Notations

Before we can start with the construction of the asymptotic solution we must introduce more definitions and notations. In particular, we must introduce parallel manifolds to the manifold \( \Gamma \) and we must extend the definition of the surface gradients for functions defined on \( \Gamma \), which are given in Section 2.1 to functions defined on the parallel manifolds. These definitions and notations are needed throughout the remaining sections.

Let \( \delta > 0 \) be the number from (2.3). For \( \xi \) satisfying \( -\delta < \xi < \delta \)

\[
\Gamma_{\xi} = \{(t, \eta + n(t, \eta)\xi) \mid (t, \eta) \in \Gamma\}
\]

is a three dimensional parallel manifold of \( \Gamma \) embedded in \( U_\delta \), and

\[
\Gamma_{\xi}(t) = \{x \in \Omega \mid (t, x) \in \Gamma_{\xi}\}
\]

is a two-dimensional parallel surface of \( \Gamma(t) \) embedded in \( U_\delta(t) \). Let \( \tau_1, \tau_2 \in \mathbb{R}^3 \) be two orthogonal unit vectors tangent to \( \Gamma_{\xi}(t) \) at \( x \in \Gamma_{\xi}(t) \). For functions \( w : \Gamma_{\xi}(t) \to \mathbb{R} \), \( W : \Gamma_{\xi}(t) \to \mathbb{R}^3 \) and \( \hat{W} : \Gamma_{\xi}(t) \to \mathbb{R}^{3 \times 3} \) we define the surface gradient and the surface divergence on \( \Gamma_{\xi}(t) \) by

\[
\begin{align*}
\nabla_{\Gamma_{\xi}} w &= (\partial_{\tau_1} w)\tau_1 + (\partial_{\tau_2} w)\tau_2, \quad (3.1) \\
\nabla_{\Gamma_{\xi}} W &= (\partial_{\tau_1} W) \otimes \tau_1 + (\partial_{\tau_2} W) \otimes \tau_2, \quad (3.2) \\
\text{div}_{\Gamma_{\xi}} W &= \tau_1 \cdot \partial_{\tau_1} W + \tau_2 \cdot \partial_{\tau_2} W = \sum_{i=1}^{2} \tau_i \cdot (\nabla_{\Gamma_{\xi}} W)\tau_i, \quad (3.3) \\
\text{div}_{\Gamma_{\xi}} \hat{W} &= (\partial_{\tau_1} \hat{W})\tau_1 + (\partial_{\tau_2} \hat{W})\tau_2. \quad (3.4)
\end{align*}
\]

Clearly, with \( \nabla_{\Gamma} \) defined in (2.5) and (2.6) we have \( \nabla_{\Gamma_{\xi}} = \nabla_{\Gamma} \). For brevity we write \( \text{div}_{\Gamma} = \text{div}_{\Gamma_0} \). If \( w, W, \hat{W} \) are defined on \( \Gamma_{\xi} \), we define \( \nabla_{\Gamma_{\xi}} w : \Gamma_{\xi} \to \mathbb{R}^3 \), \( \nabla_{\Gamma_{\xi}} W : \Gamma_{\xi} \to \mathbb{R}^{3 \times 3} \), \( \text{div}_{\Gamma_{\xi}} W : \Gamma_{\xi} \to \mathbb{R} \), \( \text{div}_{\Gamma_{\xi}} \hat{W} : \Gamma_{\xi} \to \mathbb{R} \) by applying the operators \( \nabla_{\Gamma_{\xi}} \) and \( \text{div}_{\Gamma_{\xi}} \) to the restrictions \( w|_{\Gamma_{\xi}(t)}, W|_{\Gamma_{\xi}(t)}, \hat{W}|_{\Gamma_{\xi}(t)} \) for every \( t \). With these definitions we have the splittings

\[
\begin{align*}
\nabla_x w(t, x) &= \partial_{\xi} w(t, \eta, \xi) n(t, \eta) + \nabla_{\Gamma_{\xi}} w(t, \eta, \xi), \quad (3.5) \\
\nabla_x W(t, x) &= \partial_{\xi} W(t, \eta, \xi) \otimes n(t, \eta) + \nabla_{\Gamma_{\xi}} W(t, \eta, \xi), \quad (3.6) \\
\text{div}_x \hat{W}(t, x) &= \left( (\partial_{\xi} \hat{W}(t, \eta, \xi)) n(t, \eta) + \text{div}_{\Gamma_{\xi}} \hat{W}(t, \eta, \xi) \right), \quad (3.7)
\end{align*}
\]

where, as usual, \( W(t, \eta, \xi) = W(t, \eta + n(t, \eta)\xi) \).
The operators $\nabla^\Gamma$ and $\text{div}^\Gamma$ can be applied to functions defined on subsets of $\Gamma$. In contrast, the operator $\nabla^J$ introduced next can be applied to functions defined on $\Gamma \times J$, where $J \subseteq \mathbb{R}$ is an interval. For $w : \Gamma \times J \to \mathbb{R}$, $W : \Gamma \times J \to \mathbb{R}^3$, $\hat{W} : \Gamma \times J \to \mathbb{R}^{3 \times 3}$ consider the functions $\eta \mapsto w_{t,\xi}(\eta) = w(t, \eta, \xi)$, $\eta \mapsto W_{t,\xi}(\eta) = W(t, \eta, \xi)$, $\eta \mapsto \hat{W}_{t,\xi}(\eta)$ which are defined on $\Gamma(t)$. To these functions the operators $\nabla^J$ and $\text{div}^J$ can be applied. We set

$$\nabla^\eta w(t, \eta, \xi) = \nabla^\Gamma w_{t,\xi}(\eta) \in \mathbb{R}^3,$$  

$$\nabla^\eta W(t, \eta, \xi) = \nabla^\Gamma W_{t,\xi}(\eta) \in \mathbb{R}^{3 \times 3},$$  

$$\text{div}^\eta W(t, \eta, \xi) = \text{div}^\Gamma W_{t,\xi}(\eta) \in \mathbb{R},$$  

$$\text{div}^\eta \hat{W}(t, \eta, \xi) = \text{div}^\Gamma \hat{W}_{t,\xi}(\eta) \in \mathbb{R}^3.$$  

If $W$ is defined on $\mathcal{U}_\delta$, then $(t, \eta, \xi) \to W(t, \eta, \xi) = W(t, \eta + n(t, \eta)\xi)$ is defined on $\Gamma \times (-\delta, \delta)$. Consequently, the gradient $\nabla^\eta W$ is defined. The connection between $\nabla^\eta W$ and $\nabla^\Gamma W = \nabla^\Gamma W|_{\Gamma^t}$ is given by the chain rule, which yields

$$\nabla^\eta W(t, \eta, \xi) = (\nabla^\Gamma W(t, \eta + n(t, \eta)\xi))(I + \xi \nabla^\eta n(t, \eta)).$$  

(3.12)

In particular, we have $\nabla^\eta W(t, \eta, 0) = \nabla^\Gamma W(t, \eta)$. Similar formulas and relations hold for $\nabla^\eta w$, $\text{div}^\eta W$, $\text{div}^\eta \hat{W}$. If $W : \mathcal{U}_\delta \to \mathbb{R}^3$ is constant on all the lines normal to $\Gamma(t)$, for all $t$, we have $W(t, \eta, \xi) = W(t, \eta)$. For such functions we sometimes interchangeably use the notations $\nabla^\eta W$ and $\nabla^\Gamma W$. Similarly, we interchangeably use the notations $\nabla^\eta w$ and $\nabla^\Gamma w$, $\text{div}^\eta W$ and $\text{div}^\Gamma W$, $\text{div}^\eta \hat{W}$ and $\text{div}^\Gamma \hat{W}$ if $w$ and $\hat{W}$ are independent of $\xi$.

Note that by (3.12) we have for $x \in \Gamma^t(t)$ that

$$\nabla^\Gamma W(t, x) = \left(\nabla^\eta W(t, \eta, \xi)\right)A(t, \eta, \xi),$$  

(3.13)

where $A(t, \eta, \xi) \in \mathbb{R}^{3 \times 3}$ is the inverse of the linear mapping $(I + \xi \nabla^\eta n(t, \eta)) : \mathbb{R}^3 \to \mathbb{R}^3$. From the mean value theorem we obtain the expansion

$$A(t, \eta, \xi) = I + \xi R_A(t, \eta, \xi)$$  

(3.14)

where the remainder term $R_A(t, \eta, \xi) \in \mathbb{R}^{3 \times 3}$ is bounded when $(t, \eta, \xi)$ varies in $\Gamma \times (-\delta, \delta)$. Insertion into (3.13) yields

$$\nabla^\Gamma W(t, x) = \nabla^\eta W(t, \eta, \xi) \left(I + \xi R_A(t, \eta, \xi)\right).$$  

(3.15)

For $w : \mathcal{U}_\delta \to \mathbb{R}$ we consider $\nabla^\Gamma w$ and $\nabla^\eta w$ to be column vectors. For such $w$ the equation corresponding to (3.15) is

$$\nabla^\Gamma w(t, x) = A^T(t, \eta, \xi) \nabla^\eta w(t, \eta, \xi) = \left(I + \xi R_A^T(t, \eta, \xi)\right) \nabla^\eta w(t, \eta, \xi).$$  

(3.16)

Furthermore, (3.3), (3.15) and (3.10) together yield for $W : \mathcal{U}_\delta \to \mathbb{R}^3$ that

$$\text{div}^\Gamma W = \sum_{i=1}^2 \tau_i \cdot ((\nabla^\eta W)(I + \xi R_A)\tau_i) = \text{div}^\eta W + \xi \text{div}^\Gamma \xi W,$$

(3.17)

with the remainder term

$$\text{div}^\Gamma \xi W(t, \eta, \xi) = \sum_{i=1}^2 \tau_i \cdot ((\nabla^\eta W)R_A\tau_i) = \sum_{i=1}^2 \tau_i \cdot ((\nabla^\Gamma W_{t,\xi})R_A\tau_i),$$  

(3.18)
and this equation implies for \( \hat{W} : \mathcal{U}_\delta \rightarrow \mathbb{R}^{3 \times 3} \) that

\[
\text{div}_{\Gamma_\xi} \hat{W}(t, \eta, \xi) = \text{div}_\eta \hat{W} + \xi \text{div}_{\Gamma_\xi} \hat{W}, \quad (3.19)
\]

where \( \text{div}_{\Gamma_\xi} \hat{W} = \sum_{i,j=1}^{2} (\partial_{\tau_j} \hat{W}_{i,\xi}) \tau_i (\tau_j \cdot R_{A_\tau}) \). The terms \( \text{div}_{\Gamma_\xi} W \) and \( \text{div}_{\Gamma_\xi} \hat{W} \) are bounded when \( (t, \eta, \xi) \) varies in \( \Gamma \times (-\delta, \delta) \).

For functions \( w \) with values in \( \mathbb{R} \) we define the second gradients \( \nabla_\eta^2 w, \nabla_\eta^2 \hat{W} \) by applying the operators \( \nabla_{\Gamma_\xi}^2, \nabla_\eta \) to the vector functions \( \nabla_{\Gamma_\xi} w, \nabla_\eta w \). For \( \hat{W} \) with values in \( \mathbb{R}^3 \) we define second gradients \( \nabla_{\Gamma_\xi}^2 \hat{W}, \nabla_\eta^2 \hat{W} \) by applying these operators to the rows of \( \nabla_{\Gamma_\xi} \hat{W}, \nabla_\eta \hat{W} \). We remark that

\[
\Delta_{\Gamma_\xi} w = \text{div}_{\Gamma_\xi} \nabla_{\Gamma_\xi} w
\]

is the surface Laplacian.

**Definition 3.1** Let \( I \subseteq \mathbb{R} \) be an interval. For \( k, m \in \mathbb{N}_0 \) and \( p = 1, 3 \) we define the space

\[
C_k(I, C^m(\Gamma, \mathbb{R}^p)) = \{(t, \eta, \xi) \rightarrow w(t, \eta, \xi) : \Gamma \times I \rightarrow \mathbb{R}^p \mid \partial_\xi^\ell \partial_i^j \nabla_\eta^2 w \in C(\Gamma \times I), \ell \leq k, i + j \leq m \}.
\]

### 3.2 Construction of the asymptotic solution

We start with the construction of the asymptotic solution \((u^{(\mu)}, T^{(\mu)}, S^{(\mu)})\). We assume that the hypotheses of Theorem 2.3 are satisfied. In particular, we assume that there is a sufficiently smooth solution \( \Gamma = \Gamma^{(\mu)}(u) \) of the evolution problem (2.14), (2.15), with \( s_0(T, \kappa_T, \lambda^{1/2}), s_1(\tilde{u}, \tilde{T}, S_0, S_1, \lambda^{1/2}) \) defined in (2.40) – (2.43). By this assumption, the function \((\tilde{u}, \tilde{T}, \tilde{u}, \tilde{T}, S_0, S_1)\) is known as a solution of the transmission-boundary value problem (2.16) – (2.32). We use the notation

\[
1^+(r) = \begin{cases} 
1, & r > 0 \\
0, & r \leq 0
\end{cases}, \quad 1^-(r) = 1 - 1^+(r), \quad r^\pm = r 1^\pm(r). \quad (3.20)
\]

Let \( \phi \in C^\infty_0((-2, 2)) \) be a function satisfying \( 0 \leq \phi(r) \leq 1 \) for all \( r \in \mathbb{R} \) and \( \phi(r) = 1 \) for \( |r| \leq 1 \). With the constant \( a \) from (2.33) we define a function \( \phi_{\mu \lambda} : Q \rightarrow [0, 1] \) by

\[
\phi_{\mu \lambda}(t, x) = \phi_{\mu \lambda}(t, \eta, \xi) = \phi\left(\frac{2a_\xi}{3(\mu \lambda)^{1/2} |\ln \mu|}\right), \quad \text{for } (t, x) \in \mathcal{U}_\delta, \quad (3.21)
\]

\[
\phi_{\mu \lambda}(t, x) = 0, \quad \text{otherwise.}
\]

By (2.34), \( \phi_{\mu \lambda} \) is equal to 1 in \( Q^{(\mu \lambda)}_{\text{in}} \), transits smoothly from 1 to 0 in \( Q^{(\mu \lambda)}_{\text{match}} \), and vanishes in \( Q^{(\mu \lambda)}_{\text{out}} \). With this function the asymptotic solution is defined by

\[
u^{(\mu)}(t, x) = u_1^{(\mu)}(t, x) \phi_{\mu \lambda}(t, x) + u_2^{(\mu)}(t, x) (1 - \phi_{\mu \lambda}(t, x)), \quad (3.22)
\]

\[
S^{(\mu)}(t, x) = S_1^{(\mu)}(t, x) \phi_{\mu \lambda}(t, x) + S_2^{(\mu)}(t, x) (1 - \phi_{\mu \lambda}(t, x)), \quad (3.23)
\]

\[
T^{(\mu)}(t, x) = D \left( \varepsilon (\nabla_x u^{(\mu)}(t, x)) - \tau S^{(\mu)}(t, x) \right), \quad (3.24)
\]

25
where \( u^{(\mu)}_1, S^{(\mu)}_1 \) are components of the inner expansion \( (u^{(\mu)}_1, T^{(\mu)}_1, S^{(\mu)}_1) \) defined in \( U_\delta \), and \( u^{(\mu)}_2, S^{(\mu)}_2 \) are components of the outer expansion \( (u^{(\mu)}_2, T^{(\mu)}_2, S^{(\mu)}_2) \) defined in \( Q \setminus \Gamma \). The function \( (u^{(\mu)}, T^{(\mu)}, S^{(\mu)}) \) is equal to the inner expansion \( (u^{(\mu)}_1, T^{(\mu)}_1, S^{(\mu)}_1) \) in the region \( Q^{(\mu)}_{\text{in}} \) and equal to the outer expansion \( (u^{(\mu)}_2, T^{(\mu)}_2, S^{(\mu)}_2) \) in the region \( Q^{(\mu)}_{\text{out}} \). In the region \( Q^{(\mu)}_{\text{match}} \) both expansions are matched.

**The outer expansion** The outer expansion is defined as follows. With the solutions \((\tilde{u}, \tilde{T})\) of the transmission problem (2.16) – (2.20) and \((\hat{u}, \hat{T})\) of the transmission problem (2.21) – (2.25) we set for \((t, x) \in Q \setminus \Gamma\)

\[
\begin{align*}
u^{(\mu)}_2(t, x) &= \hat{u}(t, x) + \mu^{1/2} \tilde{u}(t, x) + \mu \tilde{u}(t, x), \\
S^{(\mu)}_2(t, x) &= \hat{S}(t, x) + \mu^{1/2} \tilde{S}_1(t, x) + \mu \tilde{S}_2(t, x) + \mu^{3/2} \tilde{S}_3(t, x), \\
T^{(\mu)}_2(t, x) &= D \left( \varepsilon (\nabla u^{(\mu)}_2(t, x)) - \varepsilon S^{(\mu)}_2(t, x) \right).
\end{align*}
\]

(3.25)

(3.26)

(3.27)

The functions \( \tilde{u}, \tilde{S}_1, \ldots, \tilde{S}_3 \) and another unknown function \( \tilde{T} \) solve the system of algebraic and partial differential equations

\[
\begin{align*}
div_x \tilde{T} &= 0, \\
\tilde{T} &= D \left( \varepsilon (\nabla x \hat{u}) - \varepsilon \hat{S}_2 \right), \\
-\tilde{T} : \varepsilon + \hat{\psi}'(\hat{S}) \tilde{S}_1 &= 0, \\
-\tilde{T} : \varepsilon + \hat{\psi}'(\hat{S}) \tilde{S}_2 + \frac{1}{2} \hat{\psi}''(\hat{S}) \tilde{S}_1^2 &= 0, \\
-\tilde{T} : \varepsilon + \hat{\psi}'(\hat{S}) \tilde{S}_3 + \hat{\psi}''(\hat{S}) \tilde{S}_1 \tilde{S}_2 + \frac{1}{6} \hat{\psi}^{(IV)}(\hat{S}) \tilde{S}_1^3 + \frac{\lambda^{1/2}}{c} \partial_t \tilde{S}_1 - \lambda \Delta_x \tilde{S}_1 &= 0,
\end{align*}
\]

(3.28)

(3.29)

(3.30)

(3.31)

(3.32)

in the set \( Q \setminus \Gamma \). Moreover, \( \tilde{u} \) satisfies the boundary conditions

\[
\begin{align*}
\tilde{u}(t, x) &= 0, \quad (t, x) \in [t_1, t_2] \times \partial \Omega, \\
\tilde{u}^{(-)}(t, \eta) &= \lambda^{1/2} a^*(t, \eta) \int_{-\infty}^{0} S_1(t, \eta, \zeta) - \frac{\varepsilon \cdot \hat{T}^{(-)}(t, \eta, \zeta)}{\hat{\psi}'(1)} d\zeta, \quad (t, \eta) \in \Gamma, \\
\tilde{u}^{(+)}(t, \eta) &= \lambda^{1/2} a^*(t, \eta) \int_{0}^{\infty} S_1(t, \eta, \zeta) - \frac{\varepsilon \cdot \hat{T}^{(+)}(t, \eta, \zeta)}{\hat{\psi}'(1)} d\zeta \\
&\quad + \lambda a^*(t, \eta) \int_{-\infty}^{\infty} \left( \int_{-\infty}^{\zeta} S_0(\vartheta) d\vartheta - \zeta^+ \right) d\zeta, \quad (t, \eta) \in \Gamma.
\end{align*}
\]

(3.33)

(3.34)

(3.35)

Since by assumption \( \Gamma, \hat{T}, \hat{T}, S_1 \) are known from the evolution problem, this system can be solved recursively. To see this, note that (3.30) yields

\[
\tilde{S}_1 = \frac{\hat{T} : \varepsilon}{\hat{\psi}''(\hat{S})}.
\]

(3.36)

We insert this equation into (3.31) and solve this equation for \( \tilde{S}_2 \) to obtain

\[
\tilde{S}_2 = \frac{\hat{T} : \varepsilon}{\hat{\psi}''(\hat{S})} - \frac{\hat{\psi}''(\hat{S})}{2 \hat{\psi}''(\hat{S})} \left( \frac{\hat{T} : \varepsilon}{\hat{\psi}''(\hat{S})} \right)^2.
\]

(3.37)
Using this function in (3.29), we can determine $\tilde{u}$ and $\tilde{T}$ from the boundary value problem (3.28), (3.29), (3.33) – (3.35). Finally, we can solve (3.32) for $S_3$.

The inner expansion The inner expansion $(u_1^{(\mu)}, T_1^{(\mu)}, S_1^{(\mu)})$ is essentially obtained by smoothing the jumps of the functions $\hat{u}$, $\check{u}$ and $\hat{S}$ from the evolution problem for the surface $\Gamma$. Before we can define the inner expansion we must therefore study in the next two lemmas the jumps of $\hat{u}$ and $\check{u}$ across $\Gamma$.

Let $u^* = [\partial_n \hat{u}]$ and $a^* = [\partial_n^2 \hat{u}]$ be the jumps of derivatives of $\hat{u}$ across $\Gamma$. These functions are introduced in (2.35), (2.36). For $(t,x) = (t,x(t,\eta,\xi)) \in U_\delta$ we decompose $\hat{u}$ and $\check{u}$ in the form

\[
\hat{u}(t,x) = u^*(t,\eta) \xi^+ + a^*(t,\eta) \frac{1}{2} (\xi^+)^2 + \hat{v}(t,x),
\]

\[
\check{u}(t,x) = u^*(t,\eta) \left( \frac{\bar{\psi}^+}{\psi^+(1)} \xi^+ + \frac{\bar{\psi}^-}{\psi^-(0)} \xi^- \right) + \check{v}(t,x),
\]

where $\xi^+$, $\xi^-$ are defined in (3.20) and where the remainder terms $\hat{v}$ and $\check{v}$ are defined by (3.38), (3.39). The decomposition (3.38) is motivated by the fact that

\[
[\partial_n^i \hat{v}] = 0, \quad i = 0, 1, 2,
\]

which follows immediately from (3.35), (2.18) and (2.35), (2.36). The first two terms on the right hand side of (3.38) thus serve to separate off the jumps of the first and second derivatives of $\hat{u}$ at $\Gamma$. Similarly, the normal derivatives of first order of $\hat{v}$ do not jump across $\Gamma$. More precisely, we have the following result.

Lemma 3.2 Let $(\hat{u}, \hat{T})$ and $(\check{u}, \check{T})$ be solutions of the transmission problems (2.16) – (2.20) and (2.21) – (2.25), respectively.

(i) $\hat{v}$ defined in (3.39) satisfies

\[
[\partial_n^i \hat{v}] = 0, \quad i = 0, 1.
\]

(ii) Assume that $\Gamma$ is a $C^5$–manifold. Suppose that $\hat{u} \in C^4(\gamma \cup \gamma', \mathbb{R}^3)$, $\check{u} \in C^3(\gamma \cup \gamma', \mathbb{R}^3)$ and that $\hat{u}$ has $C^4$–extensions, $\check{u}$ has $C^3$–extensions from $\gamma$ to $\gamma \cup \Gamma$ and from $\gamma'$ to $\gamma' \cup \Gamma$.

With the function spaces introduced in Definition 3.1 we then have

\[
\hat{v} \in C^2\left((-\delta, \delta), C^2(\Gamma)\right) \cap C^3\left((-\delta, 0], C^1(\Gamma)\right) \cap C^3\left([0, \delta), C^1(\Gamma)\right),
\]

\[
\check{v} \in C^1\left((-\delta, \delta), C^2(\Gamma)\right) \cap C^3\left(\Gamma \times (-\delta, 0]\right) \cap C^3\left(\Gamma \times [0, \delta)\right).
\]

Proof: To prove (3.41) note that by definition of $[\partial_n w]$ in Section 2.1 and by definition of $\xi^\pm$ in (3.20) we have

\[
\left[\partial_n \left( \frac{\bar{\psi}^+}{\psi^+(1)} \xi^+ + \frac{\bar{\psi}^-}{\psi^-(0)} \xi^- \right) u^* \right] = \left( \frac{\bar{\psi}^+}{\psi^+(1)} - \frac{\bar{\psi}^-}{\psi^-(0)} \right) u^* = \left[ \frac{\bar{\psi}^+}{\psi^+(S)} \right] u^*.
\]

From this equation, from the first equation in (2.38) and from (3.39) we obtain (3.41) for $i = 1$. For $i = 0$ equation (3.41) is an immediate consequence of (3.39) and (2.23). This proves (i).
(ii) Since $\Gamma$ is a $C^5$–manifold, the coordinate mapping $(t, \eta, \xi) \mapsto (t, x(t, \eta, \xi)) = (t, \eta + \xi n(t, \eta))$ and the inverse mapping $(t, x) \mapsto (t, \eta(t, x), \xi(t, x))$ are $C^4$. It follows from this differentiability property of the coordinate mapping and from our differentiability assumptions for $\mathbf{u}$ that $(t, \eta, \xi) \mapsto \mathbf{u}(t, \eta, \xi)$ is $C^4$ in $\Gamma \times (-\delta, 0]$ and in $\Gamma \times [0, \delta)$, and that $(t, \eta) \mapsto \mathbf{u}^*(t, \eta) = n(t, \eta) \cdot [\nabla_x \mathbf{u}(t, \eta)]$ belongs to $C^3(\Gamma)$ and $(t, \eta) \mapsto \mathbf{a}^*(t, \eta) = n(t, \eta) \cdot [\partial_n \nabla_x \mathbf{u}(t, \eta)]$ belongs to $C^2(\Gamma)$. Since by (3.38) we have

$$\mathbf{v}(t, \eta, \xi) = \mathbf{u}(t, \eta, \xi) - \mathbf{u}^*(t, \eta) \xi^* - \mathbf{a}^*(t, \eta) \frac{1}{2} (\xi^*)^2,$$

these properties imply that

$$\mathbf{v} \in \bigcap_{m=0}^1 \left( C^{2+m}((-\delta, 0], C^{2-m}(\Gamma)) \cap C^{2+m}([0, \delta), C^{2-m}(\Gamma)) \right).$$

From this relation and from (3.40) we conclude that (3.42) holds. Relation (3.43) is obtained in the same way, using (3.41) instead of (3.40).

For brevity in notation we define

$$\mathbf{\hat{\sigma}}(\xi) = \mathbf{\hat{\sigma}}(t, \eta, \xi) = \varepsilon : \mathbf{\hat{T}}(\xi) = D\varepsilon \left( \nabla_x \mathbf{\hat{v}}(t, x) \right), \quad \mathbf{\hat{\sigma}}'(\xi) = \partial_\varepsilon \mathbf{\hat{\sigma}}(t, \eta, \xi),$$

$$\mathbf{\hat{\sigma}}(\xi) = \mathbf{\hat{\sigma}}(t, \eta, \xi) = \varepsilon : \mathbf{\hat{T}}(\xi) = D\varepsilon \left( \nabla_x \mathbf{\hat{v}}(t, x) \right).$$

(3.44)

(3.45)

Later we need the following result, which shows how $\mathbf{\hat{\sigma}}(t, \eta, 0)$ can be computed from the limit values of $\mathbf{\hat{T}}$ and $\mathbf{\hat{T}}$ at $\Gamma$.

**Lemma 3.3** The function $\mathbf{\hat{\sigma}}$ defined in (3.45) satisfies

$$\varepsilon : \mathbf{\hat{T}}(\xi) = \mathbf{\hat{\sigma}}(0) + \varepsilon : \mathbf{\hat{T}} \frac{\varepsilon : \mathbf{\hat{T}}(\xi)}{\psi''(1)},$$

(3.46)

$$\varepsilon : \mathbf{\hat{T}}(\xi) = \mathbf{\hat{\sigma}}(0) + \varepsilon : \mathbf{\hat{T}} \frac{\varepsilon : \mathbf{\hat{T}}(\xi)}{\psi''(0)},$$

(3.47)

$$\mathbf{\hat{\sigma}}(0) = \varepsilon : \mathbf{\hat{T}} - \varepsilon : \mathbf{\hat{T}} \frac{\psi''(S)}{\psi''(1)}.$$  (3.48)

**Proof:** We apply the decomposition (2.28) of the gradient to the function $W(t, \eta, \xi) = u^*(t, \eta) \left( \frac{\varepsilon : \mathbf{\hat{T}}(\xi)}{\psi''(1)} \xi^* + \frac{\varepsilon : \mathbf{\hat{T}}(\xi)}{\psi''(0)} \xi^* \right)$. This yields

$$\left( \nabla_x W \right)^{(+)}(t, \eta) = (\partial_n W) \otimes n + \nabla_\Gamma W = u^* \otimes n \frac{\varepsilon : \mathbf{\hat{T}}(\xi)}{\psi''(1)}.$$  (3.39)

(3.39) thus implies

$$\left( \nabla_x \mathbf{\hat{u}} \right)^{(+)} = u^* \otimes n \frac{\varepsilon : \mathbf{\hat{T}}(\xi)}{\psi''(1)} + \nabla_x \mathbf{\hat{v}}.$$

Insertion of these equations into (2.22) yields

$$\mathbf{\hat{T}}^{(+)} = D\varepsilon (u^* \otimes n) \frac{\varepsilon : \mathbf{\hat{T}}(\xi)}{\psi''(1)} + D\varepsilon (\nabla_x \mathbf{\hat{v}}).$$

We take the scalar product with $\varepsilon$ on both sides of this equation and note (2.95) and the definition of $\mathbf{\hat{\sigma}}$ in (3.45) to obtain (3.46). Equation (3.47) is obtained in the same way. To prove (3.48), we add (3.46) and (3.47) and solve the resulting equation for $\mathbf{\hat{\sigma}}(0)$. This proves the lemma.
**Definition of the inner expansion** We can now construct the inner expansion \((u_1^{(\mu)}, T_1^{(\mu)}, S_1^{(\mu)})\). With the remainder terms \(\hat{v}, \hat{v}\) introduced in (3.38), (3.39) we set in the neighborhood \(U_0\) of \(\Gamma\)

\[
u_1^{(\mu)}(t, x) = (\mu \lambda)^{1/2} u_0(t, \eta, \frac{\xi}{(\mu \lambda)^{1/2}}) + \mu \lambda^{1/2} u_1(t, \eta, \frac{\xi}{(\mu \lambda)^{1/2}}) + \mu \lambda u_2(t, \eta, \frac{\xi}{(\mu \lambda)^{1/2}}) + \hat{v}(t, x) + \mu^{1/2} \hat{v}(t, x),
\]

(3.49)

\[
S_1^{(\mu)}(t, x) = S_0\left(\frac{\xi}{(\mu \lambda)^{1/2}}\right) + \mu^{1/2} S_1(t, \eta, \frac{\xi}{(\mu \lambda)^{1/2}}) + \mu S_2(t, \eta, \frac{\xi}{(\mu \lambda)^{1/2}}),
\]

(3.50)

\[
T_1^{(\mu)}(t, x) = D\left(\varepsilon(\nabla_x u_1^{(\mu)}(t, x)) - \varepsilon S_1^{(\mu)}(t, x)\right),
\]

(3.51)

where by assumption \(S_0, S_1\) are known from the evolution problem for \(\Gamma\), and where the functions \(u_0, \ldots, u_2\) are defined by

\[
u_0(t, \eta, \zeta) = u^*(t, \eta) \int_{-\infty}^{\zeta} S_0(\vartheta)\,d\vartheta,
\]

(3.52)

\[
u_1(t, \eta, \zeta) = u^*(t, \eta) \int_{0}^{\zeta} S_1(t, \eta, \vartheta)\,d\vartheta,
\]

(3.53)

\[
u_2(t, \eta, \zeta) = a^*(t, \eta) \int_{-\infty}^{\zeta} \int_{-\infty}^{\vartheta} S_0(\vartheta_1)\,d\vartheta_1\,d\vartheta,
\]

(3.54)

The function \(S_2 = S_2(t, \eta, \zeta)\) together with another unknown function \(s_1 = s_1(t, \eta)\) solve a boundary value problem. To state this boundary value problem let \(\kappa(t, \eta, \xi)\) denote twice the mean curvature of the surface \(\Gamma_\xi(t)\) at \(\eta \in \Gamma_\xi(t)\). With the notation introduced in Section 2.2 we thus have \(\kappa(t, \eta, 0) = \kappa_\Gamma(t, \eta, 0)\).

The boundary value problem for \(S_2\) and \(s_1\) consists of the ordinary differential equation

\[
\dot{\psi}''(S_0(\zeta)) S_2(t, \eta, \zeta) - S_2''(t, \eta, \zeta) = F_2(t, \eta, \zeta),
\]

(3.55)

with the right hand side given by

\[
F_2(t, \eta, \zeta) = \dot{\sigma}(0) + \varpi: [\dot{T}] S_1 - \frac{1}{c_1} \varpi: (\dot{T}) S_1' - \frac{1}{2} \psi''(S_0) S_1^2
\]

\[
+ \lambda^{1/2} \left(\dot{\sigma}'(0) \zeta + \varpi: D\varepsilon(a^* \otimes n + \nabla_\Gamma u^*) \int_{-\infty}^{\zeta} S_0(\vartheta)\,d\vartheta\right)
\]

\[
+ \left(\frac{s_1}{c} - \lambda \kappa'(0) \zeta\right) S_0',
\]

(3.56)

and of boundary conditions. To formulate these boundary conditions, we choose \(\varphi \in C^\infty(\mathbb{R}, [0, 1])\) such that

\[
\varphi(\zeta) = \begin{cases} 0, & \zeta \leq 1, \\ 1, & \zeta \geq 2, \end{cases}
\]

(3.57)

set

\[
\varphi_+(\zeta) = \frac{\varphi(\zeta)}{\psi''(1)}, \quad \varphi_-(\zeta) = \frac{\varphi(-\zeta)}{\psi''(0)},
\]

(3.58)
and define

\[
\rho_2(t, \eta, \zeta) = \varphi_-(\zeta) \left( \tau : \mathcal{T}(-) - \frac{\psi'''(0)}{2} \left( \frac{\tau : \mathcal{T}(-)}{\psi''(0)} \right)^2 + \lambda^{1/2} \sigma'(0) \zeta \right) \\
+ \varphi_+(\zeta) \left( \tau : \mathcal{T}(+) - \frac{\psi'''(1)}{2} \left( \frac{\tau : \mathcal{T}(+)}{\psi''(1)} \right)^2 + \lambda^{1/2} \sigma'(0) \zeta \\
+ \lambda^{1/2} \varepsilon : (a^* \otimes n + \nabla \Gamma u^*) \zeta \right). 
\]

(3.59)

With this function the boundary conditions are

\[
S_2(t, \eta, 0) = 0, \quad \lim_{\zeta \to \pm \infty} (S_2(t, \eta, \zeta) - \rho_2(t, \eta, \zeta)) = 0. 
\]

(3.60)

The function \( s_1 = s_1(t, \eta) \) in (3.56) is independent of \( \zeta \). It is determined in Section 4.2 by the procedure sketched at the end of Section 2.2, which we apply of course to the boundary value problem (3.55), (3.56), (3.60), (3.61) instead of the problem (2.27), (2.29) – (2.32). The function \( s_1 \), whose explicit expression is given in (2.41) – (2.43), forms the second term in the definition (2.15) of the evolution operator \( K^{(\nu)} \).

4 The functions \( S_0, \ldots, S_2 \) from the inner expansion

The functions \( \tilde{S}_1, \ldots, \tilde{S}_3 \) in the outer expansion can be determined explicitly from (3.30) – (3.32), whereas the functions \( S_0, \ldots, S_2 \) in the inner expansion are determined as solutions of three coupled boundary value problems to linear and nonlinear ordinary differential equations. It is not obvious that these solutions exist and what properties they have. We study these solutions in this section.

4.1 The function \( S_0 \)

The first boundary value problem determining \( S_0 \) is given by (2.26), (2.28).

Lemma 4.1 Assume that the double well potential \( \hat{\psi} \) satisfies (2.38). Then \( S_0 \) is a solution of the boundary value problem (2.26), (2.28), if and only if \( S_0 \) satisfies the initial value problem

\[
S_0'(\zeta) = \sqrt{2\hat{\psi}(S_0(\zeta))}, \quad \zeta \in \mathbb{R}, \quad S_0(0) = \frac{1}{2}. 
\]

(4.1)

Proof: Let \( S_0 \) be a solution of (2.26), (2.28). We multiply (2.26) by \( S_0' \) and obtain

\[
\frac{d}{d\zeta} \left( \hat{\psi}(S_0) - \frac{1}{2} (S_0')^2 \right) = 0,
\]

or

\[
\hat{\psi}(S_0) - \frac{1}{2} (S_0')^2 = C_1. 
\]

(4.2)

By (2.28) we have \( \lim_{\zeta \to \pm \infty} S_0(\zeta) = 1 \). From (4.2) and from (2.38) we thus obtain that \( \lim_{\zeta \to \pm \infty} (S_0'(\zeta))^2 = -2C_1 \). Using again (2.28), we infer from this limit relation that
\[ \lim_{\zeta \to \infty} S_0'(\zeta) = 0, \text{ hence } C_1 = 0. \] We solve (4.2) for \( S_0' \) and use that because of the boundary conditions (2.28) the function \( S_0 \) must be increasing, hence \( S_0' \) must be non-negative. This shows that a solution of (2.26), (2.28) must satisfy the initial value problem (4.1).

To prove the converse we differentiate the differential equation in (4.1) and obtain (2.26). We leave it to the reader to verify that the solution of (4.1) satisfies the boundary conditions (2.28).

**Theorem 4.2** Assume that \( \hat{\psi} \in C^3([0,1], \mathbb{R}) \) has the properties (2.38). Then there is a unique solution \( S_0 \in C^4(\mathbb{R}, (0,1)) \) of the initial value problem (4.1). This solution is strictly increasing and satisfies (2.26) and (2.28). Moreover, there are constants \( K_1, \ldots, K_3 > 0 \) such that for \( a > 0 \) defined in (2.38)

\[
0 < S_0(\zeta) \leq K_1 e^{-a|\zeta|}, \quad -\infty < \zeta \leq 0, \tag{4.3}
\]

\[
1 - K_2 e^{-a\zeta} \leq S_0(\zeta) < 1, \quad 0 \leq \zeta < \infty, \tag{4.4}
\]

\[
|\partial^i S_0(\zeta)| \leq K_3 e^{-a|\zeta|}, \quad -\infty < \zeta < \infty, \quad i = 1, \ldots, 4. \tag{4.5}
\]

This theorem follows immediately from the standard theory of ordinary differential equations, and we omit the proof.

**Lemma 4.3** If \( \hat{\psi} \) satisfies the symmetry condition (2.39), then the solution \( S_0 \) of (4.1) satisfies for all \( \zeta \in \mathbb{R} \)

\[
S_0(-\zeta) = 1 - S_0(\zeta), \quad S_0'(\zeta) = S_0'(-\zeta), \tag{4.6}
\]

\[
\int_{-\infty}^\zeta S_0(\vartheta) \, d\vartheta = \int_{-\infty}^{-|\zeta|} S_0(\vartheta) \, d\vartheta + \zeta^+, \tag{4.7}
\]

\[
\left| \int_{-\infty}^\zeta S_0(\vartheta) \, d\vartheta - \zeta^+ \right| \leq \frac{K_1}{a} e^{-a|\zeta|}, \tag{4.8}
\]

\[
|\hat{\psi}''(S_0(\zeta)) - \hat{\psi}''(\hat{S}(\zeta))| \leq K_4 e^{-a|\zeta|}. \tag{4.9}
\]

**Proof:** If the symmetry condition (2.39) holds and if \( S_0 \) is a solution of the initial value problem (4.1), then also \( \zeta \mapsto (1 - S_0(-\zeta)) \) is a solution. To see this, note that (2.39) and (4.1) imply

\[
\sqrt{2\hat{\psi}}\left(1 - S_0(-\zeta)\right) = \sqrt{2\hat{\psi}\left(\frac{1}{2} + \left(\frac{1}{2} - S_0(-\zeta)\right)\right)} = \sqrt{2\hat{\psi}\left(\frac{1}{2} - (\frac{1}{2} - S_0(-\zeta))\right)}
\]

\[
= \sqrt{2\hat{\psi}(S_0(-\zeta))} = (\partial_\zeta S_0)(-\zeta) = \partial_\zeta (1 - S_0(-\zeta)),
\]

whence \( 1 - S_0(-\zeta) \) satisfies the differential equation in (4.1). Since we obviously have \( 1 - S_0(0) = \frac{1}{2} \), we see that \( 1 - S_0(-\zeta) \) is a solution of (4.1). Since the solution of this initial value problem is unique, we infer that \( S_0(\zeta) = 1 - S_0(-\zeta) \) holds, which implies \( S_0'(\zeta) = S_0'(-\zeta) \).
To prove (4.17), note that (4.16) implies for \( \zeta > 0 \)
\[
\int_{-\infty}^{\zeta} S_0(\vartheta) \, d\vartheta = \int_{-\infty}^{-\zeta} S_0(\vartheta) \, d\vartheta + \int_{-\zeta}^{\zeta} S_0(\vartheta) \, d\vartheta \\
= \int_{-\infty}^{-\zeta} S_0(\vartheta) \, d\vartheta + \int_{0}^{\zeta} S_0(\vartheta) + S_0(-\vartheta) \, d\vartheta \\
= \int_{-\infty}^{-\zeta} S_0(\vartheta) \, d\vartheta + \int_{0}^{\zeta} S_0(\vartheta) + (1 - S_0(\vartheta)) \, d\vartheta = \int_{-\infty}^{-\zeta} S_0(\vartheta) \, d\vartheta + \zeta.
\]
From this equation we immediately obtain (4.17). The inequality (4.18) follows from (4.17) and from (4.3), which yield
\[
\left| \int_{-\infty}^{\zeta} S_0(\vartheta) \, d\vartheta - \zeta^+ \right| = \left| \int_{-\infty}^{-|\zeta|} S_0(\vartheta) \, d\vartheta \right| \leq \int_{-\infty}^{-|\zeta|} K_1 e^{-a|\vartheta|} \, d\vartheta = \frac{K_1}{a} e^{-a|\zeta|}.
\]

For the proof of (4.19) note that \( \hat{S}(\zeta) = 1 \) for \( \zeta > 0 \). Consequently, the mean value theorem and (4.21) together imply for \( \zeta > 0 \) that
\[
|\hat{\psi}''(S_0(\zeta)) - \hat{\psi}''(\hat{S}(\zeta))| = |\hat{\psi}''(S_0(\zeta)) - \hat{\psi}''(1)| \leq |\hat{\psi}''((S_0(\zeta) - 1)| \leq C K_2 e^{-a|\zeta|},
\]
with a suitable number \( r^* \) between \( S_0(\zeta) \) and 1. For \( \zeta < 0 \) an analogous estimate is obtained using (4.3) and noting that \( \hat{S}(\zeta) = 0 \) if \( \zeta < 0 \).

### 4.2 The functions \( S_1 \) and \( S_2 \)

The solutions \( S_1 \) and \( S_2 \) of the second and third boundary value problems are studied in this section. The second problem determining \( S_1 \) is given by the equations (2.27), (2.29) – (2.32), the third problem, which determines \( S_2 \), consists of the equations (3.55), (3.56), (3.60), (3.61). The properties of \( S_1 \) and \( S_2 \), which we need in Section 5, are summarized in the next two theorems.

To state the first theorem we need the function \( \rho_1 : \Gamma \times \mathbb{R} \to \mathbb{R} \), which is defined by
\[
\rho_1(t, \eta, \zeta) = \overline{\zeta} : \hat{T}^{(-)}(t, \eta) \varphi_-(\zeta) + \overline{\zeta} : \hat{T}^{(+)}(t, \eta) \varphi_+(\zeta), \tag{4.10}
\]
with \( \varphi\pm \) introduced in (3.58).

**Theorem 4.4** Assume that \( \hat{\psi} \) belongs to \( C^5([0, 1], \mathbb{R}) \) and satisfies the assumptions (2.38) and the symmetry condition (2.39). Suppose that the function \( s_0 = s_0(t, \eta) \) in (2.32) is given by (2.40). Let \( S_0 \) be the solution of the boundary value problem (2.26), (2.28), which exists by Theorem 4.2.

Then for every \( (t, \eta) \in \Gamma \) there is a unique solution \( \zeta \mapsto S_1(t, \eta, \zeta) : \mathbb{R} \to \mathbb{R} \) of the boundary value problem (2.27), (2.29) – (2.32). The function \( S_1 \) belongs to the space \( C^2(\mathbb{R}, C^2(\Gamma, \mathbb{R})) \). Moreover, there are constants \( K_1, \ldots, K_3 \) such that for the constant \( a \) defined in (2.38) and for all \( (t, \eta, \zeta) \in \Gamma \times \mathbb{R} \) the estimates
\[
\|D_{(t, \eta)}^{\alpha} S_1\|_{L^\infty(\Gamma \times \mathbb{R})} \leq K_1, \quad |\alpha| \leq 2, \tag{4.11}
\]
\[
|\hat{\partial}_\zeta D_{(t, \eta)}^{\alpha} S_1(t, \eta, \zeta) - \rho_1(t, \eta, \zeta)| \leq K_2 e^{-a|\zeta|}, \quad 0 \leq j, |\alpha| \leq 2, \tag{4.12}
\]
\[
|\hat{\partial}_\zeta^j D_{(t, \eta)}^{\alpha} S_1(t, \eta, \zeta)| \leq K_3 e^{-a|\zeta|}, \quad |\alpha| \leq 2, \quad j = 1, 2, \tag{4.13}
\]
hold.
We do not give the proof of this theorem, since it is almost the same as the proof of Theorem 1.2 in [7]. Morover, it is obtained from the proof of the following theorem by simplification. The main difference between the two proofs is that the right hand side $F_1$ of the differential equation (2.27) for $S_1$ is bounded, whereas the right hand side $F_2$ of the differential equation (3.55) for $S_2$ grows linearly for $\zeta \to \pm \infty$.

**Theorem 4.5** Assume that $\hat{\psi}$ satisfies the assumptions given in the Theorem 4.4. Let $S_0$ be the solution of the boundary value problem (2.26), (2.28), and let $S_1$ be the solution of the boundary value problem (2.27), (2.29) – (2.32). Suppose that the function $s_1 = s_1(t,\eta)$ in (3.56) satisfies (2.41) with $s_{10}$, $s_{11}$ given in (2.42), (2.43).

(i) Then for every $(t,\eta) \in \Gamma$ there is a unique solution $\zeta \mapsto S_2(t,\eta,\zeta) : \mathbb{R} \to \mathbb{R}$ of the boundary value problem (3.55), (3.56), (3.60), (3.61). The function $S_2$ belongs to $C^2(\mathbb{R},C^2(\Gamma,\mathbb{R}))$, and there are constants $K_4,\ldots,K_6$ such that for the constant $a$ defined in (2.38) and for all $(t,\eta,\zeta) \in \Gamma \times \mathbb{R}$ the estimates

$$|\partial^j_\zeta D^{\alpha}_{(t,\eta,\zeta)} S_2(t,\eta,\zeta)| \leq K_4 (1 + |\zeta|)^{1-j}, \quad |\alpha| \leq 2, \quad j = 0,1,$$  

$$|\partial^j_\zeta D^{\alpha}_{(t,\eta)} (S_2(t,\eta,\zeta) - \rho_2(t,\eta,\zeta))| \leq K_5 (1 + |\zeta|) e^{-a|\zeta|}, \quad 0 \leq j, |\alpha| \leq 2,$$  

$$|\partial^2_\zeta^2 D^{\alpha}_{(t,\eta)} S_2(t,\eta,\zeta)| \leq K_6 (1 + |\zeta|) e^{-a|\zeta|}, \quad |\alpha| \leq 2,$$

hold, where $\rho_2$ is defined in (3.59).

(ii) $S_2$ is the only solution of the differential equation (3.55) with $F_2$ given by (3.56), which satisfies (3.60) and for which constants $C,\theta > 0$ exist such that

$$|S_2(t,\eta,\zeta)| \leq C e^{(\alpha-\theta)|\zeta|}, \quad \zeta \in \mathbb{R},$$

holds.

**4.3 Proof of Theorem 4.5**

In this section we give the proof of Theorem 4.5, which is divided into five parts:

(I) **Reduction of the boundary value problem for $S_2$ to a problem in $L^2$.** With $\rho_2$ defined in (3.59) we make the ansatz

$$S_2(t,\eta,\zeta) = w(t,\eta,\zeta) + \rho_2(t,\eta,\zeta).$$  

Insertion of this ansatz into the equations (3.55) and (3.60), (3.61) shows that $S_2$ is a solution of the problem given by these equations if and only if $w$ solves the equations

$$\hat{\psi}''(s_0(\zeta)) w(t,\eta,\zeta) - \partial^2_\zeta w(t,\eta,\zeta) = F_2(t,\eta,\zeta) + F_3(t,\eta,\zeta),$$

$$w(t,\eta,0) = 0,$$  

$$\lim_{\zeta \to \pm \infty} w(t,\eta,\zeta) = 0,$$

where $F_2$ is given by (3.56) and where

$$F_3 = - (\hat{\psi}''(s_0) - \partial^2_\zeta) \rho_2.$$  

To get (4.20) we used that $\varphi_+(0) = \varphi_-(0) = 0$, which by (3.59) implies $\rho_2(t,\eta,0) = 0$. To show that the solution $S_2$ of the problem (3.55) and (3.60), (3.61) exists, it therefore suffices to prove that the reduced problem (4.19) – (4.22) has a solution.
(II) Spectral theory  For this proof note that \( \hat{\psi}''(S_0) - \partial_\zeta^2 \) is a linear self-adjoint differential operator in \( L^2(\mathbb{R}) \). From the spectral theory of such operators we know that the continuous spectrum of \( \hat{\psi}''(S_0) - \partial_\zeta^2 \) is contained in the interval \( [a_0, \infty) \), where

\[
a_0 = \min \left\{ \lim_{\zeta \to -\infty} \hat{\psi}''(S_0(\zeta)), \lim_{\zeta \to \infty} \hat{\psi}''(S_0(\zeta)) \right\},
\]

and that the part of the spectrum in \( (-\infty, a_0) \) is a pure point spectrum. \((4.9)\) yields

\[
\lim_{\zeta \to -\infty} \hat{\psi}''(S_0(\zeta)) = \hat{\psi}''(0), \quad \lim_{\zeta \to \infty} \hat{\psi}''(S_0(\zeta)) = \hat{\psi}''(1),
\]

hence the assumption \((2.38)\) implies \( a_0 = a^2 > 0 \). Therefore 0 does not belong to the continuous spectrum. From the spectral theory we also know that for every \( \omega \in \mathbb{C} \), which is not in the continuous spectrum, the differential equation \((\hat{\psi}''(S_0) - \partial_\zeta^2)w - \omega w = f\) has a solution \( w \in L^2(\mathbb{R}) \), if and only if \( f \in L^2(\mathbb{R}) \) is orthogonal to the kernel of the operator \( \hat{\psi}''(S_0) - \partial_\zeta^2 - \omega \). This implies in particular, that for every \( (t, \eta) \in \Gamma \) the differential equation \((4.19)\) has a solution \( w(t, \eta, \zeta) \in L^2(\mathbb{R}) \), if the right hand side \( \zeta \mapsto f(\zeta) = F_2(t, \eta, \zeta) + F_3(t, \eta, \zeta) \) belongs to \( L^2(\mathbb{R}) \) and is orthogonal to the kernel of the operator \( \hat{\psi}''(S_0) - \partial_\zeta^2 \). To show that the problem \((4.19) - (4.21)\) has a solution, we therefore verify in the next two parts of the proof that \( F_2 + F_3 \) satisfies these two conditions.

(III) The asymptotic behavior of \( F_2 + F_3 \) at infinity.  We first show that the right hand side \( F_2 + F_3 \) of \((4.19)\) decays exponentially at \( \pm \infty \), which implies that \( F_2 + F_3 \in L^2(\mathbb{R}) \). To simplify the notation we define

\[
\varphi_+(\zeta) = \varphi_+(\zeta)\hat{\psi}''(S_0(\zeta)), \quad \varphi_- (\zeta) = \varphi_- (\zeta)\hat{\psi}''(S_0(\zeta)),
\]

with \( \varphi_+, \varphi_- \) given in \((3.58)\). For these functions we obtain from \((4.9)\) that

\[
|\varphi_- - \hat{\psi}''(0)\varphi_-| = |\varphi_- (\zeta) |\hat{\psi}''(S_0) - \hat{\psi}''(0)| \leq CK_4 e^{-a|\zeta|}, \quad (4.24)
\]

\[
|\varphi_+ - \hat{\psi}''(1)\varphi_+| = |\varphi_+ (\zeta) |\hat{\psi}''(S_0) - \hat{\psi}''(1)| \leq CK_4 e^{-a|\zeta|}, \quad (4.25)
\]

for all \( \zeta \in \mathbb{R} \). Since \( \hat{\psi}''(0)\varphi_- (\zeta) = 1 \) for \( \zeta \leq -2 \) and \( \hat{\psi}''(1)\varphi_+ (\zeta) = 1 \) for \( \zeta \geq 2 \), these estimates imply

\[
|1 - \varphi_- (\zeta)| \leq CK_4 e^{-a|\zeta|}, \quad -\infty < \zeta \leq 0. \quad (4.26)
\]

\[
|1 - \varphi_+ (\zeta)| \leq CK_4 e^{-a|\zeta|}, \quad 0 \leq \zeta < \infty, \quad (4.27)
\]

\[
|1 - \varphi_- (\zeta) - \varphi_+ (\zeta)| \leq CK_4 e^{-a|\zeta|}, \quad \zeta \in \mathbb{R}. \quad (4.28)
\]

To get the last estimate we combined the first two estimates and noted that \( \varphi_- (\zeta) = 0 \) for \( \zeta \geq -1 \) and \( \varphi_+ (\zeta) = 0 \) for \( \zeta \leq 1 \).

Note that by \((3.56)\), \((3.59)\) and \((4.22)\) the function \( F_2 + F_3 \) can be decomposed in the form

\[
F_2 + F_3 = F_2 - \hat{\psi}''(S_0)\rho_2 + \partial_\zeta^2 \rho_2 = \sum_{j=1}^{5} I_j, \quad (4.29)
\]
where

\[
I_1 = \bar{\sigma}(0) + \bar{\tau} : [\bar{T}] S_1 - \bar{\tau} : \hat{T}^{(-)} \varphi_+ - \bar{\tau} : \hat{T}^{(+)} \varphi_+,
\]

\[
I_2 = -\frac{\hat{\psi}''''(S_0)}{2} S_1^2 + \frac{\hat{\psi}''''(0)}{2} \left( \bar{\tau} : \hat{T}^{(-)} \right) \varphi_- - \frac{\hat{\psi}''''(1)}{2} \left( \bar{\tau} : \hat{T}^{(+)} \right) \varphi_+,
\]

\[
I_3 = \lambda^{1/2} \bar{\sigma}'(0) \zeta (1 - \varphi_+ - \varphi_-),
\]

\[
I_4 = \lambda^{1/2} \bar{\tau} : D \varepsilon (a^* \otimes n + \nabla_x u^*) \left( \int_{-\infty}^{\zeta} S_0 D \vartheta \, d\vartheta - \zeta^+ \varphi_+ \right),
\]

\[
I_5 = -\frac{1}{c_1} \bar{\tau} : \langle \hat{T} \rangle S'_1 + \left( \frac{s_1}{c} - \lambda \kappa'(0) \zeta \right) S'_0 + \partial_2 \rho_2.
\]

We show that everyone of these terms decays to zero for \( \zeta \to \pm \infty \). To verify this for the first term we insert \((3.46)\) and \((3.47)\) into \((4.30)\), which results in

\[
|I_1| \leq |\bar{\sigma}(0) (1 - \varphi_+ + \varphi_-) + \bar{\tau} : \hat{T} \left( S_1 - \rho_1 \right)|
\]

\[
+ \left| \bar{\tau} : \hat{T} \left( \bar{\tau} : \hat{T}^{(-)} \varphi_- - \hat{\psi}''''(0) \varphi_- \right) + \bar{\tau} : \hat{T}^{(+)} \varphi_+ - \hat{\psi}''''(1) \varphi_+ \right| \leq Ce^{-a|\zeta|},
\]

for all \( \zeta \in \mathbb{R} \), where we applied the estimates \((4.31)\), \((4.34)\), \((4.35)\) and \((4.38)\).

Next we estimate \( I_2 \). By definition we have \( \varphi_+(\zeta) = 0 \) for \( \zeta \leq 1 \). From \((4.31)\) we thus have on the half axis \( -\infty < \zeta \leq 0 \) that

\[
I_2 = -\frac{\hat{\psi}''''(S_0)}{2} S_1^2 + \frac{\hat{\psi}''''(0)}{2} \left( \bar{\tau} : \hat{T}^{(-)} \right) \varphi_- - \frac{\hat{\psi}''''(1)}{2} \left( \bar{\tau} : \hat{T}^{(+)} \right) \varphi_+.
\]

To estimate \( I_2 \) we apply the mean value theorem to \( \hat{\psi}'''' \) and use \((4.3)\) and \((4.11)\), to estimate \( I_{22} \) we use \((4.26)\) and \((4.11)\). The result is

\[
|I_2| \leq CK_1 e^{-a|\zeta|}, \quad -\infty < \zeta \leq 0.
\]

To estimate \( I_{23} \) note that by \((3.57)\), \((3.58)\) and \((4.10)\) we have for \( -\infty < \zeta \leq -2 \) that

\[
\rho_1(t, \eta, \zeta) = \varphi_-(\zeta) \bar{\tau} : \hat{T}^{(-)}(t, \eta) = \frac{\bar{\tau} : \hat{T}^{(-)}(t, \eta)}{\psi''''(0)}.
\]

With this equation we infer from \((4.11)\) and \((4.12)\) with \( \alpha, \beta = 0 \) that

\[
|I_{23}| \leq \left| \frac{\hat{\psi}''''(0)}{2} \left( S_1 - \bar{\tau} : \hat{T}^{(-)} \right) \left( S_1 + \bar{\tau} : \hat{T}^{(-)} \right) \varphi_- \right| \leq Ce^{-a|\zeta|}, \quad -\infty < \zeta \leq 0.
\]
\[ (4.36) - (4.38) \] together imply that \(|I_2(\zeta)| \leq C e^{-a|\zeta|} \) for \(-\infty < \zeta \leq 0\). On the half axis \(0 \leq \zeta < \infty\) we estimate \(I_2\) analogously. This proves that

\[ |I_2| \leq C e^{-a|\zeta|}, \quad -\infty < \zeta < \infty. \quad (4.39) \]

The estimate for \(I_3\) is obtained by application of \((4.28)\) to \((4.32)\), which immediately yields

\[ |I_3| \leq C (1 + |\zeta|) e^{-a|\zeta|}, \quad -\infty < \zeta < \infty. \quad (4.40) \]

To study the asymptotic behavior of \(I_4\) note that \((4.8)\) and \((4.27)\) together imply

\[ \int_{-\infty}^{\zeta} S_0(\vartheta) d\vartheta - \zeta^+ \overline{\varphi}(\zeta) \leq \left| \int_{-\infty}^{\zeta} S_0(\vartheta) d\vartheta - \zeta^+ \right| + \zeta^+ |1 - \overline{\varphi}(\zeta)| \leq \left( \frac{1}{a} K_1 + \zeta^+ C K_4 \right) e^{-a|\zeta|}. \]

Insertion of this inequality into \((4.33)\) results in

\[ |I_4| \leq C (1 + |\zeta|) e^{-a|\zeta|}. \quad (4.41) \]

It remains to investigate \(I_5\). Note first that the third term on the right hand side of \((4.34)\) satisfies

\[ \partial_\xi^2 \rho(t, \eta, \zeta) = 0, \quad \text{for } |\zeta| \geq 2. \quad (4.42) \]

To show this it suffices to remark that the functions \(\varphi_{\pm}\) are constant on the intervals \((-\infty, -2)\) and \((2, \infty)\), from which we see by inspection of \((3.59)\) that on these intervals the function \(\zeta \mapsto \rho_2(t, \eta, \zeta)\) is a sum of constant and linear terms, whence \((4.42)\) follows.

If we estimate the first term on the right hand side of \((4.34)\) by employing \((4.13)\) with \(\alpha = 0, j = 1\) and the second term by using \((4.5)\), we obtain together with \((4.42)\) that

\[ |I_5| \leq C (1 + |\zeta|) e^{-a|\zeta|}. \quad (4.43) \]

We combine \((4.29), (4.35), (4.39) - (4.41)\) and \((4.43)\) to derive the estimate

\[ |F_2(t, \eta, \zeta) + F_3(t, \eta, \zeta)| \leq C (1 + |\zeta|) e^{-a|\zeta|}, \quad \text{for all } \zeta \in \mathbb{R}, \quad (4.44) \]

which shows in particular that the right hand side of \((4.19)\) belongs to \(L^2(\mathbb{R})\).

\textbf{(IV) The orthogonality condition determining } s_1. \quad \text{Next we must show that the right hand side of } (4.19) \text{ is orthogonal to the kernel of } \hat{\psi}''(S_0) - \partial_\xi^2. \text{ This kernel is different from } \{0\}, \text{ since } S_0' \text{ belongs to the kernel. This is immediately seen by differentiation of } (2.26), \text{ which yields}

\[ \hat{\psi}''(S_0(\zeta)) S_0'(\vartheta) - \partial_\xi^2 S_0'(\zeta) = 0. \quad (4.45) \]

Since by \((4.5)\) the function \(S_0'\) is in the domain of definition of \(\hat{\psi}''(S_0) - \partial_\xi^2\), it belongs to the kernel of this operator.

The theory of linear ordinary differential equations of second order implies now that the kernel is one-dimensional, hence every function from the kernel is a multiple of \(S_0'\). Therefore the right hand side \(F_2 + F_3\) of \((4.19)\) is orthogonal to the kernel if it is orthogonal to \(S_0'\). Note that the integrals \(\int_{-\infty}^{\infty} F_2 S_0' d\zeta\) and \(\int_{-\infty}^{\infty} F_3 S_0' d\zeta\) both exist, since \(F_2\) and \(F_3\)
grow at most linearly for $\zeta \to \pm \infty$, whereas by (4.5) the function $S'_0$ decays exponentially at $\pm \infty$. Therefore we obtain from (1.22) and (4.5) by partial integration that

$$\int_{-\infty}^{\infty} (F_2 + F_3)S'_0 d\zeta = \int_{-\infty}^{\infty} F_2S'_0 d\zeta - \int_{-\infty}^{\infty} \left( (\hat{\psi}''(S_0) - \partial_\zeta^2) \rho_2 \right) S'_0 d\zeta$$

$$= \int_{-\infty}^{\infty} F_2S'_0 d\zeta - \int_{-\infty}^{\infty} \rho_2(\hat{\psi}''(S_0) - \partial_\zeta^2) S'_0 d\zeta = \int_{-\infty}^{\infty} F_2S'_0 d\zeta. \quad (4.46)$$

To study the last integral on the right hand side note that by (4.6) the function $S'_0$ is even, which implies that

$$\int_{-\infty}^{\infty} \hat{\sigma}'(t, \eta, 0) \zeta S'_0(\zeta) d\zeta = 0, \quad \text{and} \quad \int_{-\infty}^{\infty} \kappa'(t, \eta, 0) \zeta (S'_0(\zeta))^2 d\zeta = 0. \quad (4.47)$$

Moreover, since by Lemma 4.1 the function $S'_0$ satisfies (4.1), we obtain by substitution of $\theta = S_0(\zeta)$

$$\int_{-\infty}^{\infty} S'_0(\zeta) S'_0(\zeta) d\zeta = \int_{-\infty}^{1} \sqrt{2\hat{\psi}'(S_0(\zeta))} S'_0(\zeta) d\zeta = \int_{0}^{1} \sqrt{2\hat{\psi}'(\theta)} d\theta = c_1, \quad (4.48)$$

where the last equality sign holds by definition of $c_1$ in (2.37). Finally, by partial integration,

$$\int_{-\infty}^{\infty} \int_{-\infty}^{\zeta} S_0(\theta) d\theta S'_0(\zeta) d\zeta = \lim_{\zeta_1 \to \infty} \left( \int_{-\infty}^{\zeta_1} S_0(\theta) d\theta S_0(\zeta_1) - \int_{-\infty}^{\zeta_1} S_0(\zeta)^2 d\zeta \right)$$

$$= \lim_{\zeta_1 \to \infty} \int_{-\infty}^{\zeta_1} \left( S_0(\zeta_1) - S_0(\zeta) \right) d\zeta$$

$$= \int_{-\infty}^{\infty} S_0(\zeta)(1 - S_0(\zeta)) d\zeta = \int_{-\infty}^{\infty} S_0(\zeta)S_0(-\zeta) d\zeta. \quad (4.49)$$

In the second last step we used that $S_0$ is increasing. The equality sign thus follows from the theorem of Beppo Levi. The last equality sign is obtained from (4.6).

The equations (3.56) and (4.37) – (4.39) yield

$$\int_{-\infty}^{\infty} F_2(t, \eta, \zeta) S'_0(\zeta) d\zeta = \int_{-\infty}^{\infty} \left( \hat{\sigma}(0) + \tau : [\hat{T}] S_1 \right) S'_0(\zeta) d\zeta - \frac{1}{c_1} \tau : \langle \hat{T} \rangle \int_{-\infty}^{\infty} S'_1 S_0 d\zeta$$

$$- \frac{1}{2} \int_{-\infty}^{\infty} \hat{\psi}''(S_0) S_1^2 S'_0 d\zeta + \lambda^{1/2} \tau : D\varepsilon(a^* \otimes n + \nabla u^*) \int_{-\infty}^{\infty} \int_{-\infty}^{\zeta} S_0(\theta) d\theta S'_0(\zeta) d\zeta$$

$$+ \int_{-\infty}^{\infty} \left( \lambda^{1/2} \hat{\sigma}'(0) + \left( \frac{s_1}{c} - \lambda \kappa'(0) \right) S'_0 \right) S'_0(\zeta) d\zeta$$

$$= \hat{\sigma}(0) + \tau : [\hat{T}] \int_{-\infty}^{\infty} S_1 S'_0 d\zeta - \frac{1}{c_1} \tau : \langle \hat{T} \rangle \int_{-\infty}^{\infty} S'_1 S_0 d\zeta - \frac{1}{2} \int_{-\infty}^{\infty} \hat{\psi}''(S_0) S_1^2 S'_0 d\zeta$$

$$+ \lambda^{1/2} \tau : D\varepsilon(a^* \otimes n + \nabla u^*) \int_{-\infty}^{\infty} S_0(\zeta)S_0(-\zeta) d\zeta + c_1 \frac{s_1}{c}.$$

To get the second last equality sign we inserted (3.48) for $\hat{\sigma}(0)$ and used (2.42), (2.43). The right hand side of this equation vanishes if and only if $s_1$ satisfies (2.41). From (4.46)
we thus infer that $F_2 + F_3$ is orthogonal to the kernel of $\hat{\psi}''(S_0) - \partial_\xi^2$ if and only if (2.41) holds.

Consequently, from part (II) of the proof we conclude that the differential equation (4.19) has a solution $w$ in $L^2(\mathbb{R})$ if and only if $s_1$ satisfies (2.41) with $s_{10}$ and $s_{11}$ given in (2.42), (2.43). In fact, there is exactly one such $w$, which also satisfies (4.20). To prove this assume that $\tilde{w} \in L^2(\mathbb{R})$ is a special solution of (4.19). Then we obtain every solution contained in $L^2(\mathbb{R})$ in the form $w = \tilde{w} + \beta S'_0$ with an arbitrary constant $\beta \in \mathbb{R}$. Since (4.1) and (2.38) yield

$$S'_0(0) = \sqrt{2\hat{\psi}(S_0(0))} = \sqrt{2\hat{\psi}(1)} > 0,$$

we can choose $\beta$ such that

$$w(t, \eta, 0) = \tilde{w}(t, \eta, 0) + \beta S'_0(0) = 0,$$

which is (4.20). This equation determines $\beta$ uniquely, hence $w$ is the unique solution of (4.19) and (4.20) in $L^2(\mathbb{R})$.

(V) Existence of the solution, estimates (4.18) – (4.20). We show next that this function $w$ satisfies (4.21). To this end we need the following

**Lemma 4.6** Let $\hat{a}_- > 0$, $\hat{a}_+ > 0$ and set $\hat{a} = \sqrt{\min\{\hat{a}_-, \hat{a}_+\}}$. Let $g : \mathbb{R} \to \mathbb{R}$ be a smooth function and let $f : \mathbb{R} \to \mathbb{R}$ be a continuous function such that

\[
|g(\zeta) - \hat{a}_-| \leq Ce^{-\hat{a}|\zeta|}, \quad \text{for } \zeta < 0,\tag{4.50}
\]

\[
|g(\zeta) - \hat{a}_+| \leq Ce^{-\hat{a}|\zeta|}, \quad \text{for } \zeta > 0,\tag{4.51}
\]

\[
|f(\zeta)| \leq C(1 + |\zeta|)e^{-\hat{a}|\zeta|}, \quad \text{for } \zeta \in \mathbb{R}.
\tag{4.52}
\]

Let $\hat{w}$ be a solution of

$$g(\zeta)\hat{w}(\zeta) - \partial^2\hat{w}(\zeta) = f(\zeta), \quad \zeta \in \mathbb{R}.\tag{4.53}$$

(i) Then $\hat{w}$ belongs to the space $C^2(\mathbb{R})$. If $\hat{w} \in L^2(\mathbb{R})$, then there is $C > 0$ such that

$$|\partial^j\hat{w}(\zeta)| \leq C(1 + |\zeta|)e^{-\hat{a}|\zeta|}, \quad \text{for all } \zeta \in \mathbb{R}, \text{ for } j = 0, 1, 2.\tag{4.54}
\]

(ii) If there are $C, \theta > 0$ such that

$$|\hat{w}(\zeta)| \leq Ce^{(\hat{a} - \theta)|\zeta|}\tag{4.55}$$

holds for all $\zeta \in \mathbb{R}$, then $\hat{w} \in L^2(\mathbb{R})$.

This is a standard result from the theory of ordinary differential equations, and we omit the proof.

To show that $w$ satisfies (4.21), we apply this lemma with $\hat{a}_- = \hat{\psi}''(0)$, $\hat{a}_+ = \hat{\psi}''(1)$, $g(\zeta) = \hat{\psi}''(S_0(\zeta))$ and $f(\zeta) = F_2(t, \eta, \zeta) + F_3(t, \eta, \zeta)$. Then we have $\hat{a} = \sqrt{\min\{\hat{\psi}''(0), \hat{\psi}''(1)\}} = a$, by (2.38), and from (4.9) and (4.44) we see that (4.50) – (4.52) hold for this choice of functions and constants. Moreover, with this choice of functions the differential equation (4.53) is equal to (4.19). Since $w \in L^2(\mathbb{R})$ is a solution of
we see that all assumptions for part (i) of Lemma 4.6 are satisfied, hence (4.54) holds for \( w \), which means that
\[
|\partial_z^2 w(\zeta)| \leq C(1 + |\zeta|)e^{-a|\zeta|}, \quad \zeta \in \mathbb{R}, \quad j = 0, 1, 2, \tag{4.56}
\]
and this in particular implies that \( w \) satisfies (4.21).

We have now found a unique solution \( w \in L^2(\mathbb{R}) \) of (4.19) – (4.21). By part (I) of this proof this means that \( S_2 \) given by (4.18) is a solution of the boundary value problem (3.55), (3.56), (3.60) and (3.61). Since by (4.18) we have \( w = S_2 - \rho_2 \), the inequality (4.56) shows that \( S_2 \) satisfies (4.15) for \( \alpha = 0 \).

To verify that \( S_2 \) satisfies (4.15) for \( \alpha \neq 0 \), it must first be shown that \( S_2 \) is two times continuously differentiable with respect to \( (t, \eta) \). This follows if we can show that \( w = S_2 - \rho_2 \) is two times continuously differentiable with respect to \( (t, \eta) \), since by our regularity assumptions the function \( \rho_2 \) has this differentiability property. To prove this differentiability of \( w \), we write (4.19), (4.20) as a perturbation problem for the linear equation
\[
Aw = f(t, \eta)
\]
in \( L^2(\mathbb{R}) \), where \( A = (\dot{\psi}''(S_0) - \partial_z^2) \) is the linear differential operator on the left hand side of (4.19) and \( f(t, \eta) = F_2(t, \eta, \cdot) + F_3(t, \eta, \cdot) \in L^2(\mathbb{R}) \) is the function on the right hand side of (4.19), which depends two times continuously differentiable on \( (t, \eta) \) and satisfies the estimate (4.44) for every \( (t, \eta) \in \Gamma \). Since 0 is an eigenvalue of \( A \) and since \( f(t, \eta) \) is orthogonal to the kernel of \( A \) for every \( (t, \eta) \), this linear equation has infinitely many solutions and the solution set is affine. The condition (4.20) defines a linear subspace, which is closed in the Sobolev space \( H^1(\mathbb{R}) \), and which intersects the solution set in exactly one point \( w \), which is the solution of (4.19), (4.20).

To the problem set in this way we can apply the perturbation theory of linear operators. The theory yields that \( w \) is two times continuously differentiable with respect to \( (t, \eta) \). We avoid the details but refer to standard texts on the perturbation theory of linear operators, for example [26].

With this knowledge we can derive the estimate (4.15) for \( \alpha \neq 0 \) by applying the differential operator \( D_{(t,\eta)}^\alpha \) with \( 1 \leq |\alpha| \leq 2 \) to the differential equation (4.19) and obtain
\[
\dot{\psi}''(S_0)(D_{(t,\eta)}^\alpha w) - \partial_z^2(D_{(t,\eta)}^\alpha w) = D_{(t,\eta)}^\alpha (F_2 + F_3). \tag{4.57}
\]
This is a differential equation for the function \( D_{(t,\eta)}^\alpha w \) with right hand side satisfying the estimate
\[
|D_{(t,\eta)}^\alpha (F_2 + F_3)| \leq C(1 + |\zeta|)e^{-a|\zeta|}. \tag{4.58}
\]
The proof of this estimate proceeds in the same way as the proof of the corresponding estimate for \( \alpha = 0 \), which we gave in part (III). Essentially one has to replace the terms appearing in \( F_2 \) and \( F_3 \), which depend on \( (t, \eta) \), by their derivatives. To avoid repetition of many technical details, we omit this proof.

The differential equation (4.57) has the same form as the differential equation (4.19). From (4.58) we see that the assumption (4.52) holds, hence we can apply Lemma 4.6 (i) to this differential equation, from which we see that \( w \) belongs to the space \( C^2(\mathbb{R}, C^2(\Gamma, \mathbb{R})) \) and that the inequality (4.54) holds with \( \dot{w} \) replaced by \( D_{(t,\eta)}^\alpha w \), whence we have
\[
|\partial_z^j D_{(t,\eta)}^\alpha w(\zeta)| \leq C(1 + |\zeta|)e^{-a|\zeta|}, \quad \zeta \in \mathbb{R}, \quad 0 \leq j \leq 2, \quad 1 \leq |\alpha| \leq 2.
\]
Since \( w = S_2 - \rho_2 \), this is inequality (4.15) with \( \alpha \neq 0 \). Therefore we proved that (4.15) holds for all \( 0 \leq |\alpha| \leq 2 \).

The inequality (4.16) is a consequence of (4.15) and of (4.12), the inequality (4.14) follows by combination of (4.15) with the estimate

\[
|\partial_2^j D_{(t,\eta)}^\alpha \rho_2(t, \eta, \zeta)| \leq C(1 + |\zeta|)^{1-j}, \quad j = 0, 1,
\]

which is seen to hold by inspection of (3.59).

We have now proved statement (i) of Theorem 4.5 and it remains to verify (ii). That is, we have to show that \( S_2 \) is the only solution of (3.55), (3.56), (3.60) satisfying (4.17). Indeed, from (4.14) it follows that \( S_2 \) satisfies (4.17). Assume that \( S_2^* \) is a second solution satisfying (4.17). Then \( \hat{w} = S_2 - S_2^* \) fulfills (4.55) and the equation

\[
\hat{w}''(S_0) \hat{w} - \partial^2 \hat{w} = 0.
\]

Lemma 4.6(ii) thus yields \( \hat{w} \in L^2(\mathbb{R}) \). Consequently, by (4.18) we have \( S_2^* = w + \hat{w} + \rho_2 \), where \( w + \hat{w} \in L^2(\mathbb{R}) \) is a solution of (4.19), (4.20). At the end of part (IV) of this proof we showed that \( w \) is the only solution of (4.19), (4.20) in \( L^2(\mathbb{R}) \), whence \( \hat{w} = 0 \), hence \( S_2^* = S_2 \).

The proof of Theorem 4.5 is complete.

5 Proof of the estimates (2.51) – (2.54) in Theorem 2.3

The proof of (2.51) – (2.54) is straightforward: We insert the function \( (u^{(\mu)}, T^{(\mu)}, S^{(\mu)}) \) defined in Section 3 into the model equations (1.11) and (1.3) and compute the residues. However, the necessary computations are long. Therefore we divide them into four parts:

In Section 5.1 we compute for the functions \( T_1^{(\mu)} \) and \( T_2^{(\mu)} \) from the inner and outer expansions of \( T^{(\mu)} \) the residues \( \text{div}_\Gamma T_1^{(\mu)} + b \) and \( \text{div}_\Gamma T_2^{(\mu)} + b \) separately. Likewise, in Section 5.2 we insert the inner expansion \( (u_1^{(\mu)}, T_1^{(\mu)}, S_1^{(\mu)}) \) and the outer expansion \( (u_2^{(\mu)}, T_2^{(\mu)}, S_2^{(\mu)}) \) into (1.3) and compute and estimate the residues separately. With these residues we can prove (2.51) – (2.54) in the regions \( Q_{\text{inn}}^{(\mu\lambda)} \) and \( Q_{\text{out}}^{(\mu\lambda)} \), but in the matching region \( Q_{\text{match}}^{(\mu\lambda)} \) we need auxiliary estimates, which are stated in Section 5.3. All the estimates are put together in Section 5.4 to complete the proof.

5.1 Asymptotic expansion of \( \text{div}_\Gamma T^{(\mu)} + b \)

Lemma 5.1 Let \( (\hat{u}, \hat{T}) \) be the solution of the transmission problem (2.16) – (2.20). With the splitting (3.38) of \( \hat{u} \) the stress tensor field \( \hat{T} \) satisfies in the neighborhood \( \mathcal{U}_\delta \) of \( \Gamma \)

\[
\hat{T} = [\hat{T}] \hat{S} + D\varepsilon(\nabla_x \hat{v}) + D\varepsilon(a^* \otimes n + \nabla_{\Gamma_1} a^*)\xi^* + D\varepsilon(\nabla_{\Gamma_1} a^*)\frac{1}{2}(\xi^*)^2.
\]  

(5.1)

With \( \hat{\sigma}' \) defined in (3.41) and with a remainder term \( R_{\gamma^x,T} \in L^\infty(\mathcal{U}_\delta) \) we have for \( (t, \eta, \xi) \in \mathcal{U}_\delta \)

\[
\varepsilon : \hat{T}(t, \eta, \xi) = \varepsilon : \hat{T}^{(+)}(t, \eta) + \hat{\sigma}''(t, \eta, 0) \xi
\]

\[
+ \varepsilon : D\varepsilon(a^*(t, \eta) \otimes n(t, \eta) + \nabla_{\Gamma_1} a^*(t, \eta))\xi + R_{\gamma^x,T}(t, \eta, \xi)\xi^2, \quad \xi > 0,
\]

(5.2)

\[
\varepsilon : \hat{T}(t, \eta, \xi) = \varepsilon : \hat{T}^{(-)}(t, \eta) + \hat{\sigma}''(t, \eta, 0) \xi + R_{\gamma^x,T}(t, \eta, \xi)\xi^2, \quad \xi < 0.
\]  

(5.3)
The functions $u^*$ and $a^*$ introduced in (2.35), (2.36) and the normal vector $n$ satisfy on the interface $\Gamma$

\[
\begin{align*}
(D(\varepsilon(u^* \otimes n) - \sigma))n &= 0, \\
(D\varepsilon(a^* \otimes n + \nabla \Gamma u^*))n + \text{div}_\Gamma D\varepsilon(u^* \otimes n) &= 0.
\end{align*}
\]

(5.4)  (5.5)

**Proof:** With the splitting (3.6) of the gradient operator we compute from (3.38) that

\[
\begin{align*}
\nabla_x \hat{u}(t, x) &= \partial_\xi \hat{u} \otimes n + \nabla \Gamma \hat{u} \\
&= (u^*(t, \eta)1^+(\xi) + a^*(t, \eta)\xi^+) \otimes n(t, \eta) + \nabla \Gamma u^*(t, \eta)\xi^+ \\
&+ \nabla \Gamma u^*(t, \eta)\frac{1}{2}(\xi^+)^2 + \nabla_x \hat{v}(t, x).
\end{align*}
\]

We insert this equation into (2.17), note that $\hat{S}(t, x) = 1^+(\xi)$, and employ that by (2.85)

\[
D(\varepsilon(u^* \otimes n) - \sigma)\hat{S} = [\hat{T}][\hat{S}]
\]

to obtain (5.1).

By (3.15) we have

\[
\begin{align*}
\nabla \Gamma \hat{u}^*(t, \eta)\xi^+ &= \nabla \Gamma u^*(t, \eta)\xi^+ + (\nabla \Gamma u^*(t, \eta))R_A(t, \eta, \xi)(\xi^+)^2,
\end{align*}
\]

(5.6)

where we used our convention to identify $\nabla \Gamma u^*$ and $\nabla \eta u^*$, since $u^*$ does not depend on $\xi$. Noting the definition of $\hat{\sigma}$ in (3.44), we obtain from (5.1) and from (5.6) that

\[
\begin{align*}
\varepsilon : \hat{\sigma} = \varepsilon : [\hat{T}](t, \eta, \xi)^+ = \hat{\sigma}(t, \eta, \xi) \\
&+ \varepsilon : D\varepsilon(\nabla \Gamma u^*(t, \eta) + \nabla \Gamma u^*(t, \eta)R_A(t, \eta, \xi)\xi^+)^2.
\end{align*}
\]

(5.7)

(3.38) and (2.17) together imply for $\xi < 0$ that

\[
\hat{T}(t, \eta, \xi) = D\varepsilon(\nabla_x \hat{v}(t, \eta, \xi)),
\]

(5.8)

whence

\[
\varepsilon : \hat{T}(t, \eta, \xi) = \varepsilon : D\varepsilon(\nabla_x \hat{v}(t, \eta, \xi)) = \hat{\sigma}(t, \eta, \xi), \quad \xi < 0,
\]

and therefore

\[
\varepsilon : \hat{T}^-(t, \eta) = \hat{\sigma}(t, \eta, 0), \quad \varepsilon : \hat{T}^+(t, \eta) = \varepsilon : [\hat{T}](t, \eta) + \hat{\sigma}(t, \eta, 0).
\]

(5.9)

By Taylor’s formula we can express $\hat{\sigma}$ in the form

\[
\hat{\sigma}(t, \eta, \xi) = \hat{\sigma}(t, \eta, 0) + \partial_\xi \hat{\sigma}(t, \eta, 0)\xi + \partial_\xi^2 \hat{\sigma}(t, \eta, 0)|\xi|^2.
\]

We expand $\hat{\sigma}$ in (5.7) with this formula and note the equations (5.9) to obtain (5.2) and (5.3).
is an immediate consequence of (2.85) and (2.19). To prove (5.5) we apply (3.7) to calculate from (5.1) that
\[ 0 = \text{div}_x \hat{T} + b \]
\[ = \partial_\xi (\hat{T}|n \bar{S}) + \text{div}_\xi [\hat{T}] \bar{S} \]
\[ + \partial_\xi (D\varepsilon (a^* \otimes n + \nabla_{\xi} u^*) \xi^+) n + \partial_\xi \left( \frac{1}{2} D\varepsilon (\nabla_{\xi} a^*) (\xi^+)^2 \right) n \]
\[ + \text{div}_\xi D\varepsilon \left( (a^* \otimes n + \nabla_{\xi} u^*) \xi^+ + \frac{1}{2} (\nabla_{\xi} a^*) (\xi^+)^2 \right) \]
\[ + \text{div}_x D\varepsilon (\nabla_x \hat{\nu}) + b. \]  
(5.10)

From this equation we obtain for \( \xi < 0 \) that
\[ \text{div}_x D\varepsilon (\nabla_x \hat{\nu}) + b = \text{div}_x \hat{T} + b = 0. \]  
(5.11)

By assumption in Theorem 2.3 the function \( b \) is continuous at \( \Gamma \). Moreover, by Lemma 3.2 and the differentiability properties of \( \hat{\nu} \) required in Theorem 2.3 the function \( \hat{\nu} \) is two times continuously differentiable at \( \Gamma \). Therefore we infer from (5.11) that
\[ (\text{div}_x D\varepsilon (\nabla_x \hat{\nu}) + b)^{(+)} = (\text{div}_x D\varepsilon (\nabla_x \hat{\nu}) + b)^{(-)} = 0, \text{ on } \Gamma. \]  
(5.12)

With this equation and with \( \hat{T}|n = 0 \), by (2.19), we conclude from (5.10) that
\[ 0 = \lim_{\xi \to 0^+} (\text{div}_x \hat{T} + b) = (D\varepsilon (a^* \otimes n + \nabla_{\Gamma} u^*)) n + \text{div}_\Gamma [\hat{T}] \]

From this relation and from \( \text{div}_\Gamma [\hat{T}] = \text{div}_\Gamma D\varepsilon (u^* \otimes n) \), which is a consequence of (2.85), we obtain (5.5). 

Next we study the stress field \( T_1^{(\mu)} \) in the inner expansion.

**Lemma 5.2** Let \( u_1^{(\mu)} \), \( S_1^{(\mu)} \), \( T_1^{(\mu)} \) be given in (3.49) – (3.51), let \( u_0, u_1, u_2 \) be defined in (3.52) – (3.54), and let \( R_A \) be the remainder term from (3.14). We set \( \xi = \frac{\xi}{(\mu\lambda)^{1/2}} \). Then we have for \( (t, \eta, \xi) \) from the neighborhood \( U_\delta \) of \( \Gamma \)
\[ T_1^{(\mu)} (t, \eta, \xi) = \hat{T} (S_0 + \mu^{1/2} S_1) + D\varepsilon (\nabla_x (\hat{\nu} + \mu^{1/2} \hat{\nu})) \]
\[ + (\mu\lambda)^{1/2} D\varepsilon (a^* \otimes n + \nabla_{\Gamma} u^*) \int_{-\infty}^{\xi} S_0(\theta) \, d\theta + \mu R_{T_1} (\lambda, t, \eta, \xi, \zeta), \]  
(5.13)

where
\[ R_{T_1} (\lambda, t, \eta, \xi, \zeta) = D \left( \varepsilon (\nabla_{\xi} (\lambda^{1/2} u_1 + \lambda u_2) + \lambda \xi (\nabla_{\eta} u_0) R_A) - \xi S_2 \right). \]  
(5.14)

The argument of \( \hat{T} \), \( u^*, a^*, n \) is \((t, \eta)\), the argument of \( S_1, S_2, u_0, u_1, u_2 \) is \((t, \eta, \xi)\), the argument of \( \nabla_x \hat{\nu}, \nabla_x \hat{\nu}, R_A \) is \((t, \eta, \xi)\), and the argument of \( S_0 \) outside of the integral is \( (\mu\lambda)^{1/2} \). Moreover, we have
\[ \text{div}_x T_1^{(\mu)} + b = \text{div}_x D\varepsilon (\nabla_x \hat{\nu}) + b + \xi \text{div}_\xi [\hat{T}] S_0 + \mu^{1/2} \text{div}_\Gamma (D\varepsilon (\hat{T} S_0) + \mu^{1/2} \text{div}_\Gamma \xi D\varepsilon (a^* \otimes n + \nabla_{\Gamma} u^*) \int_{-\infty}^{\xi} S_0(\theta) \, d\theta \]
\[ + \mu \text{div}_x R_{T_1}. \]  
(5.15)
With \( \dot{\sigma}(0), \dot{\sigma}'(0), \dot{\sigma}(0) \) defined in (3.44), (3.45) we have

\[
\partial S \mathcal{W}(\varepsilon(\nabla_x u_1^{(\mu)}), S_1^{(\mu)})
= -\varepsilon : T_1^{(\mu)}(t, \eta, \xi) = -\varepsilon : [\dot{T}](S_0 + \mu^{1/2} S_1) - \dot{\sigma}(0) - (\mu \lambda)^{1/2} \dot{\sigma}'(0) \dot{\zeta} - \mu^{1/2} \dot{\sigma}(0)
- (\mu \lambda)^{1/2} \varepsilon : D\varepsilon(a^* \otimes n + \nabla \Gamma u^*) \int_{-\infty}^{\xi} S_0(\theta)d\theta - \mu R_W(\lambda, t, \eta, \xi, \zeta),
\]
where

\[
R_W(\lambda, t, \eta, \xi, \zeta) = \varepsilon : R_T(\lambda, t, \eta, \xi, \zeta) + \lambda \partial_\xi^2 \dot{\sigma}(t, \eta, \dot{\xi}) \dot{\zeta}^2 + \lambda^{1/2} \partial_\xi \dot{\sigma}(t, \eta, \dot{\xi}) \zeta,
\]
with suitable \( \dot{\zeta}, \dot{\xi} \) between 0 and \( \xi \) and with \( R_T \) defined in (5.14).

**Proof:** With the splitting (3.6) of the gradient operator we obtain by definition of \( u_1^{(\mu)} \) in (3.49), (3.52) – (3.54) that

\[
\nabla_x u_1^{(\mu)} = (\mu \lambda)^{1/2} \nabla_x u_0 + \mu \lambda^{1/2} \nabla_x u_1 + \mu \lambda \nabla_x u_2 + \nabla_x (\dot{v} + \mu^{1/2} \dot{v})
= (u^* \otimes n)S_0 + \mu^{1/2} u^* \otimes nS_1 + (\mu \lambda)^{1/2} a^* \otimes n \int_{-\infty}^{\xi} S_0(\theta)d\theta
+ (\mu \lambda)^{1/2} \nabla \Gamma_x u_0 + \mu \lambda^{1/2} \nabla \Gamma_x u_1 + \mu \lambda \nabla \Gamma_x u_2 + \nabla_x (\dot{v} + \mu^{1/2} \dot{v}).
\]

(3.52) and (3.15) together yield

\[
\nabla \Gamma_x u_0 = \nabla_\eta u_0 (I + \xi R_A) = \nabla \Gamma u^* \int_{-\infty}^{\xi} S_0(\theta)d\theta + (\mu \lambda)^{1/2} \zeta (\nabla_\eta u_0) R_A.
\]

We insert (5.18), (5.19) and (3.50) into (3.51). From the resulting equation we obtain (5.13) and (5.14) if we also note that by (2.85)

\[
D(\varepsilon(u^* \otimes n)(S_0 + \mu^{1/2} S_1) - \varepsilon(S_0 + \mu^{1/2} S_1))pl_1^{(\mu)} = D(\varepsilon(u^* \otimes n) - \varepsilon)(S_0 + \mu^{1/2} S_1) = [\dot{T}](S_0 + \mu^{1/2} S_1).
\]

To prove (5.15) we employ the splitting (3.7) of the divergence operator and (3.19) to compute from (5.13)

\[
\nabla_x T_1^{(\mu)} + b = \partial_\xi ([\dot{T}]nS_0) + (\text{div}_T[\dot{T}])S_0 + \xi (\text{div}_T[\dot{T}]S_0) + \mu^{1/2} (\partial_\xi ([\dot{T}]nS_1) + \text{div}_T[\dot{T}]S_1) + \mu^{1/2} \text{div}_x D\varepsilon(\nabla_x \dot{v}) + b + \mu^{1/2} \text{div}_x D\varepsilon(\nabla_x v) + (D\varepsilon(a^* \otimes n + \nabla \Gamma_x u^*) \text{div}_x R_T,
\]

By (2.85) and (5.5) we have

\[
(D\varepsilon(a^* \otimes n + \nabla \Gamma u^*) \text{div}_x R_T)S_0 + \text{div}_T[\dot{T}]S_0 = \left( (D\varepsilon(a^* \otimes n + \nabla \Gamma u^*) \text{div}_x R_T(u^* \otimes n) \right) S_0 = 0.
\]

With this equation and with (2.19) we obtain (5.15) from (5.20).
Because of this differentiability property we can apply the mean value theorem to
\(\xi < 0\), which implies \(\partial_S W(\varepsilon(\nabla x u_1^{(\mu)}), S_1^{(\mu)}) = -\tau: T_1^{(\mu)}\), and from (5.13), (5.14), using the Taylor expansions
\[
\hat{\sigma}(\xi) = \hat{\sigma}(0) + \partial_\xi \hat{\sigma}(\xi) \xi = \hat{\sigma}(0) + (\mu \lambda)^{1/2} \partial_\xi \hat{\sigma}(\xi) \xi, \\
\hat{\sigma}(\xi) = \hat{\sigma}(0) + \partial_\xi \hat{\sigma}(0) \xi + \partial_\xi^2 \hat{\sigma}(\xi) \xi^2 = \hat{\sigma}(0) + (\mu \lambda)^{1/2} \hat{\sigma}'(0) \xi + \mu \lambda \sigma''(\hat{\xi}) \xi^2.
\]
This completes the proof of Lemma 5.2.

**Corollary 5.3** Let \(Q_{\text{inn}}^{(\mu \lambda)}\) and \(Q_{\text{match}}^{(\mu \lambda)}\) be defined in (2.34) and let \(T_1^{(\mu)}\) be given by (3.51). Then there is a constant \(C\) such that for all \(0 < \mu \leq \mu_0\) and all \(0 < \lambda \leq \lambda_0\)
\[
\| \text{div}_x T_1^{(\mu)} + b \|_{L^\infty(Q_{\text{inn}}^{(\mu \lambda)} \cup Q_{\text{match}}^{(\mu \lambda)})} \leq C \left( \frac{\mu}{\lambda} \right)^{1/2} |\ln \mu|.
\]

**Proof:** We estimate the terms on the right hand side of (5.15). Note first that if \((t, \eta, \xi) \in Q_{\text{inn}}^{(\mu \lambda)} \cup Q_{\text{match}}^{(\mu \lambda)}\) and \(\zeta = \frac{\xi}{(\mu \lambda)^{1/2}},\) then
\[
|\xi| \leq \frac{3}{a} (\mu \lambda)^{1/2} |\ln \mu|, \quad |\zeta| \leq \frac{3}{a} |\ln \mu|.
\]
This follows from (2.34). With these inequalities the first two terms on the right hand side of (5.15) can be estimated as follows: From the differentiability properties of \(\hat{v}\) and \(b\), which in Theorem 2.3 are assumed to hold, it follows by Lemma 3.2 that
\[
\text{div}_x D\varepsilon(\nabla x \hat{v}) + b \in C(U_\delta) \cap C^1((-\delta, 0), C(\Gamma)) \cap C^1([0, \delta), C(\Gamma)).
\]
Because of this differentiability property we can apply the mean value theorem to \(\text{div}_x D\varepsilon(\nabla x \hat{v}) + b\), which together with (5.12) and (5.22) yields for all \((t, \eta, \xi) \in Q_{\text{inn}}^{(\mu \lambda)} \cup Q_{\text{match}}^{(\mu \lambda)}\) with \(\xi \geq 0\) that
\[
\left| \left( \text{div}_x D\varepsilon(\nabla x \hat{v}) + b \right)(t, \eta, \xi) \right| = \left| \left( \text{div}_x D\varepsilon(\nabla x \hat{v}) + b \right)^{(+)} + \partial_\xi \left( \text{div}_x D\varepsilon(\nabla x \hat{v}) + b \right)(t, \eta, \xi^*) \right| \leq C_1 \xi \leq C_2 \left( \frac{3}{a} (\mu \lambda)^{1/2} |\ln \mu| \right),
\]
with a suitable number \(\xi^*\) between 0 and \(\xi\). Since by (5.11) the term \(\text{div}_x D\varepsilon(\nabla x \hat{v}) + b\) vanishes for \(\xi < 0\), the inequality (5.24) holds for all \((t, \eta, \xi) \in Q_{\text{inn}}^{(\mu \lambda)} \cup Q_{\text{match}}^{(\mu \lambda)}\).

To estimate the last term in (5.15) note that (4.18), (4.11) and (5.22) yield
\[
0 \leq \int_{-\infty}^\xi S_0(\eta) d\eta \leq \xi^+ + C_2 \leq C_3 \left( \frac{3}{a} |\ln \mu| \right),
\]
\[
0 \leq \int_{-\infty}^\xi \int_{-\infty}^\eta S_0(\eta_1) d\eta_1 d\eta \leq \frac{1}{2} (\zeta^+)^2 + C_4 \leq C_5 \left( \frac{3}{a} |\ln \mu| \right)^2,
\]
\[
\left| \int_0^\xi S_1(t, \eta, \xi^+) d\eta \right| \leq C_6 |\xi| \leq C_6 \left( \frac{3}{a} |\ln \mu| \right).
\]
Using these inequalities, the definitions of \(u_0, u_1, u_2\) in (3.52) - (3.54) and the inequality (4.14) we obtain from (5.14) that
\[
|R_{T_1}| \leq (C_7 + \frac{C_8}{a^2}) |\ln \mu|^2,
\]
\[
|\mu \text{div}_x R_{T_1}| \leq \mu^{1/2} \lambda^{-1/2} (C_9 + \frac{C_{10}}{a^2}) |\ln \mu|^2.
\]
As introduced previously, for \( \partial_t \) slight abuse of notation could lead to confusion when we consider derivatives with respect to \( w \) which implies (5.21).

Together with (5.24) and (5.27) we find for \((t, \eta, \xi) \in Q_{\text{inn}}^{(\mu)} \cup Q_{\text{match}}^{(\mu)}\) that

\[
|\left( \text{div}_x T_1^{(\mu)} + b \right)(t, \eta, \xi)| \leq \left( \frac{C_{11}}{a^2} + C_{12} \right) \left( \mu^{1/2} + (\mu \lambda)^{1/2} |\ln \mu| + \mu^{1/2} \lambda^{-1/2} |\ln \mu|^2 \right),
\]

which implies (5.21).

In the next lemma we study the outer expansion \( T_2^{(\mu)} \).

**Lemma 5.4** Let \((\hat{u}, \hat{T})\) be the solution of the boundary value problem, which consists of the elliptic system (3.28), (3.29) with \( \hat{S}_2 \) given by (3.33), and of the boundary conditions (3.33) – (3.35). Let \( \tilde{S}_3 \) be the solution of (3.32). Then \( T_2^{(\mu)} \) defined in (3.27) satisfies on \( Q \setminus \Gamma \)

\[
T_2^{(\mu)} = \hat{T} + \mu^{1/2} \tilde{T} + \mu \tilde{T} - \mu^{3/2} D\tilde{\psi}_3, \quad (5.28)
\]

\[
\text{div}_x T_2^{(\mu)} + b = -\mu^{3/2} \text{div}_x (D\tilde{\psi}_3). \quad (5.29)
\]

**Proof:** Insertion of (3.25) and (3.26) into (3.27) yields

\[
T_2^{(\mu)} = D(\varepsilon(\nabla_x \hat{u}) - \tilde{\psi}_1) + \mu^{1/2} D(\varepsilon(\nabla_x \hat{u}) - \tilde{\psi}_1) + \mu D(\varepsilon(\nabla_x \hat{u}) - \tilde{\psi}_2) - \mu^{3/2} D\tilde{\psi}_3.
\]

Using (3.36), we see from this equation and from (2.17), (2.22), (3.29) that (5.28) holds. (5.29) is an immediate consequence of (5.28) and (2.16), (2.21), (3.28).

### 5.2 Asymptotic expansion of \( S_t + c(W_S + \hat{\psi} - \Delta_x S) \)

In this section we compute the form of the residue

\[
(\mu \lambda)^{1/2} \partial_t S + c \left( \partial_S W(\varepsilon(\nabla_x u), S) + \frac{1}{\mu^{1/2}} \hat{\psi}'(S) - \mu^{1/2} \lambda \Delta_x S \right),
\]

which is obtained when we either insert for \((u, S)\) the inner expansion \((u_1^{(\mu)}, S_1^{(\mu)})\) or the outer expansion \((u_2^{(\mu)}, S_2^{(\mu)})\) of the asymptotic solution \((u^{(\mu)}, S^{(\mu)})\).

For functions \((t, x) \mapsto w(t, x)\) defined in a neighborhood of \( \Gamma \) we write \( w(t, \eta, \xi) = w(t, x) \) with \( x = \eta + n(t, \eta) \xi \), as always. However, in the following computations this slight abuse of notation could lead to confusion when we consider derivatives with respect to \( t \). To avoid this, we introduce the notations

\[
w_i(t, x) = w_i(t, \eta, \xi) = \partial_r w(r, \eta, \xi)|_{r=t}, \quad (\partial_t w)(t, \eta, \xi) = \partial_t w(t, x).
\]

As introduced previously, for \( i = 0, 1, 2 \) we write \( S_i'(t, \eta, \xi) = \partial_{\xi_i} S_i(t, \eta, \xi) \) and \( S_i''(t, \eta, \xi) = \partial_{\xi_i}^2 S_i(t, \eta, \xi) \).
Inner expansion  We first compute (5.30) for \((u, S) = (u_1^{(\mu)}, S_1^{(\mu)})\). To this end we need

Lemma 5.5  Let \(s^{(\mu)}(t, \eta)\) be the normal speed of the phase interface \(\Gamma(t)\) at \(\eta \in \Gamma(t)\), let \(\nabla_\eta\) be the operator defined in (3.8), and let \(w\) be a function defined in a neighborhood of \(\Gamma\). Then we have

\[
\partial_t w(t, x) = w|_t(t, \eta, \xi) - \xi(\partial_t n)(t, \eta) \cdot \nabla_\eta w(t, \eta, \xi) - s^{(\mu)}(t, \eta) \partial_\xi w(t, \eta, \xi).
\]

Proof: By definition, \(\nabla_\eta w(t, \eta, \xi)\) is a tangential vector to \(\Gamma(t)\). Lemma 2.2 thus yields

\[
\partial_t w(t, x) = \partial_t w(t, \eta, \xi) = w|_t(t, \eta, \xi) + \partial_t \eta \cdot \nabla_\eta w(t, \eta, \xi) + \partial_\xi w(t, \eta, \xi) \partial_\xi
\]

\[
= w|_t(t, \eta, \xi) - \xi(\partial_t n)(t, \eta) \cdot \nabla_\eta w(t, \eta, \xi) - s^{(\mu)}(t, \eta) \partial_\xi w(t, \eta, \xi).
\]

This proves the lemma.

We apply this lemma to the function \(S_1^{(\mu)}\) defined in (3.50) to obtain

\[
\partial_t S_1^{(\mu)}(t, x) = S^{(\mu)}_{1|t}(t, \eta, \xi) - \xi(\partial_t n)(t, \eta) \cdot \nabla_\eta S_1^{(\mu)}(t, \eta, \xi) - s^{(\mu)}(t, \eta) \partial_\xi S_1^{(\mu)}(t, \eta, \xi).
\]

From this equation and from the asymptotic expansion (2.44) of \(s^{(\mu)}\) we conclude for the first term in (5.30) that

\[
(\mu \lambda)^{1/2} \partial_t S_1^{(\mu)}(t, x)
\]

\[
= (\mu \lambda)^{1/2} \partial_t \left(S_0\left(\frac{\xi}{(\mu \lambda)^{1/2}}\right) + \mu^{1/2} S_1(t, \eta, \xi, \partial_t n, \nabla_\eta S_1) + \mu R_{S_1}\right)
\]

\[
= -(s_0 + \mu^{1/2}s_1)(S_0 + \mu^{1/2}s_1' + \mu S_0') + \mu \lambda^{1/2} \tilde{R}_{S_1}
\]

\[
= -s_0 S_0' - \mu^{1/2}(s_1 S_0' + s_0 S_1') + \mu R_{S_1},
\]

with

\[
\tilde{R}_{S_1}(\mu, \lambda, t, \eta, \xi) = \left(S_1|_{t} + \mu^{1/2} S_2|_{t}\right) - \xi(\partial_t n) \cdot \nabla_\eta (S_1 + \mu^{1/2} S_2),
\]

\[
R_{S_1}(\mu, \lambda, t, \eta, \xi) = \lambda^{1/2} \tilde{R}_{S_1} - s_1 S_1' - (s_0 + \mu^{1/2}s_1) S_1'.
\]

For the third term in (5.30) we get from Taylor’s formula and from (3.50)

\[
\frac{1}{\mu^{1/2}} \frac{\partial}{\partial t} \dot{\psi}(S_0) = \frac{1}{\mu^{1/2}} \dot{\psi}'(S_0) + \dot{\psi}''(S_0) S_1 + \mu^{1/2} \frac{1}{2} \dot{\psi}'''(S_0) S_1^2 + \mu R_{\dot{\psi}},
\]

where

\[
R_{\dot{\psi}} = \frac{1}{2} \dot{\psi}'''(S_0)(2S_1 S_2 + \mu^{1/2} S_2^2) + \frac{1}{6} \dot{\psi}^{(IV)}(S_0 + \vartheta(\mu^{1/2}S_1 + \mu S_2))(S_1 + \mu^{1/2} S_2)^2.
\]

with suitable \(0 < \vartheta < 1\). Observe next that

\[
\Delta_x S_1^{(\mu)}(t, x) = \partial_\xi \nu S_1^{(\mu)}(t, \eta, \xi) - \kappa(t, \eta, \xi) \partial_\xi S_1^{(\mu)}(t, \eta, \xi) + \Delta_{\nu} S_1^{(\mu)}(t, \eta, \xi),
\]
where $\Delta_{\Gamma_\xi} = \text{div}_{\Gamma_\xi} \nabla_{\Gamma_\xi}$ denotes the surface Laplacian and where $\kappa(t, \eta, \xi)$ is twice the mean curvature of the surface $\Gamma_\xi(t)$ at $\eta \in \Gamma_\xi(t)$. With the notation $\kappa'(t, \eta, 0) = \partial_\xi \kappa(t, \eta, 0)$ we obtain from Taylor’s formula

$$
\kappa(t, \eta, \xi) = \kappa(t, \eta, 0) + \partial_\xi \kappa(t, \eta, 0) \xi + \frac{1}{2} \partial^2_\xi \kappa(t, \eta, \xi^*) \xi^2,
$$

where $R_\kappa(t, \eta, \xi) = \frac{1}{2} \partial^2_\xi \kappa(t, \eta, \xi^*)$ is the remainder term, with suitable $\xi^*$ between 0 and $\xi$. We insert (3.50) and (5.36) into (5.35) and obtain for the fourth term in (5.30) that

$$
\frac{1}{2} \lambda \Delta x S_1^{(\mu)} = \mu^{-1/2} S''_0 + S'' + \mu^{1/2} S''
$$

$$
- \lambda^{1/2} \left( \kappa_\Gamma + (\mu \lambda)^{1/2} \kappa' \xi + \mu \lambda R_\kappa^2 \right) (S'_0 + \mu^{1/2} S'_1) - \mu \lambda^{1/2} \kappa_\xi S'_2 + \mu^{1/2} \lambda \Delta x S_1^{(\mu)}
$$

$$
= \mu^{-1/2} S''_0 + \left( S''_1 - \lambda^{1/2} \kappa_\xi S'_0 \right) + \mu^{1/2} \left( S''_2 - \lambda^{1/2} \kappa_\xi S'_1 - \lambda^{1/2} \xi S''_0 \right) + \mu R_\Delta,
$$

where

$$
R_\Delta = -\lambda \kappa' \xi S'_1 - \lambda^{1/2} \kappa_\xi S'_2 + \lambda \Delta x \xi (S_1 + \mu^{1/2} S_2) + \lambda^{3/2} R_\kappa \xi^2 (S''_0 + \mu^{1/2} S'_1).
$$

From (5.31), (5.16), (5.33) and (5.37) we obtain

$$
(\mu \lambda)^{1/2} \partial_\xi S_1^{(\mu)} + c \left( \partial_\xi \mathcal{W} \mathcal{E}(\nabla_x u^{(\mu)}_1, S_1^{(\mu)}) + \frac{1}{\mu^{1/2}} \psi' (S_1^{(\mu)}) - \mu^{1/2} \lambda \Delta x S_1^{(\mu)} \right)
$$

$$
= \frac{c}{\mu^{1/2}} \left( \psi'(S_0) - S''_0 \right) + c \left( \psi''(S_0) S_1 - S''_1 - \sigma : [\hat{T}] S_0 - \hat{\sigma}(0) + \left( \lambda^{1/2} \kappa_\Gamma - \frac{8 \kappa}{c} \right) S_0 \right)
$$

$$
+ c \mu^{1/2} \left( \psi''(S_0) S_2 - S''_2 - \sigma : [\hat{T}] S_1 - \hat{\sigma}(0) + \left( \lambda^{1/2} \kappa_\Gamma - \frac{8 \kappa}{c} \right) S_1 \right)
$$

$$
- \lambda^{1/2} \left( \hat{\sigma}'(0) \xi + \hat{\sigma} : D \hat{\mathcal{E}}(\hat{a}^\ast \otimes n + \hat{\nabla}_\Gamma u^\ast) \right) \int_{-\infty}^{\zeta} S_0(\theta) d\theta
$$

$$
+ \frac{1}{2} \psi''(S_0) S_1^2 + \left( \lambda \kappa' \xi - \frac{8 \kappa}{c} \right) S_0 \right) + \mu R_{St+c(\mu)} \right),
$$

where

$$
R_{St+c(\mu)} = R_{St} + c(-R_{W} + R_{\hat{\psi}'} - R_\Delta),
$$

with $R_{St}$, $R_{W}$, $R_{\hat{\psi}'}$ and $R_\Delta$ given in (5.31), (5.17), (5.31) and (5.38), respectively.

**Corollary 5.6** Let $s_{0}$ be given by (2.40) and assume that the functions $S_0$, $S_1$ and $S_2$ satisfy the ordinary differential equations (2.26), (2.27), (3.55), with $F_1$, $F_2$ given by (2.32), (4.56). Assume moreover that the conditions (2.28) – (2.31) and (3.60), (3.61) hold. Then there is a constant $K$ such that the interior expansion $(u^{(\mu)}_1, S_1^{(\mu)} , T_1^{(\mu)})$ defined in (3.49) – (3.51) satisfies for all $(t, \eta, \xi)$ $\in Q_{\text{inn}}^{(\mu)} \cup Q_{\text{match}}^{(\mu)}$ and all $0 < \mu \leq \mu_0$, $0 < \lambda \leq \lambda_0$ the inequality

$$
\left| \partial_\xi S_1^{(\mu)} + \frac{c}{(\mu \lambda)^{1/2}} \left( \partial_\xi \mathcal{W} \mathcal{E}(\nabla_x u^{(\mu)}_1, S_1^{(\mu)}) + \frac{1}{\mu^{1/2}} \psi' (S_1^{(\mu)}) - \mu^{1/2} \lambda \Delta x S_1^{(\mu)} \right) \right|
$$

$$
= \left( \frac{\mu}{\lambda} \right)^{1/2} \left| R_{St} + c(-R_{W} + R_{\hat{\psi}'} - R_\Delta) \right| \leq K \left( \frac{\mu}{\lambda} \right)^{1/2} |\ln \mu|^2.
$$

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Proof: From (2.40) we obtain
\[ \lambda^{1/2} \kappa_{T} - \frac{s_{0}}{c} = \frac{1}{c_{1}} \tau : (\hat{T}), \]
and by (5.9) we have \( \tau : \hat{T}^{(-)} = \sigma(0) \). After insertion of these two equations into (5.39), the latter equation takes the form
\[
(\mu \lambda)^{1/2} \partial_{S} S_{1}^{(\mu)} + c \left( \partial_{S} W(\varepsilon(\nabla_{x} u_{1}^{(\mu)}), S_{1}^{(\mu)}) \right) + \frac{1}{\mu^{1/2}} \hat{\psi}'(S_{1}^{(\mu)}) - \mu^{1/2} \lambda \Delta_{x} S_{1}^{(\mu)} 
= \frac{c}{\mu^{1/2}} \left( \hat{\psi}'(S_{0} - S_{0}'') + c \left( \hat{\psi}''(S_{0}) S_{1} - S''_{1} - F_{1} \right) + c \mu^{1/2} \left( \hat{\psi}''(S_{0}) S_{2} - S''_{2} - F_{2} \right) 
+ \mu R_{S_{1} + c(\ldots)} = \mu \left( R_{S_{1}} + c(-R_{W} + R_{\psi} - R_{\Delta}) \right). \]

Here we also used (2.26), (2.27) and (3.55). Noting the inequalities (5.22) for \( \xi \) and \( \zeta \), the inequalities (4.13) - (4.15), (4.11), (4.13), (4.14) for \( S_{0}, S_{1}, S_{2} \), and the inequality (5.26) for the term \( R_{T_{1}} \), which appears in \( R_{W} \), we see by inspection of every term in (3.32), (5.17), (5.34) and (5.38) that the inequality
\[
|R_{S_{1}} + c(-R_{W} + R_{\psi} - R_{\Delta})| \leq K \ln |\mu|^{2} \]
holds. To obtain inequality (5.41) we divide (5.42) by \( (\mu \lambda)^{1/2} \) and estimate the right hand side of the resulting equation using (5.43).

Outer expansion Next we compute (5.30) for \( (u, S) = (u^{(\mu)}, S^{(\mu)}) \). Note first that (1.8) and (5.25) yield
\[
\partial_{S} W(\varepsilon(\nabla_{x} u_{2}^{(\mu)}), S_{2}^{(\mu)}) = -\tau : (\hat{T} + \mu^{1/2} \hat{T} + \mu \hat{T}) + \mu^{3/2} \tau : D \hat{x} \hat{S}_{3}. \]

Also, Taylor’s formula and (3.26) yield for a suitable \( 0 < \theta(t, x) < 1 \)
\[
\hat{\psi}'(S_{2}^{(\mu)}) = \hat{\psi}'(\hat{S}) + \hat{\psi}''(\hat{S})(\mu^{1/2} \hat{S}_{1} + \mu \hat{S}_{2} + \mu^{3/2} \hat{S}_{3})
+ \frac{1}{2} \hat{\psi}'''(\hat{S})(\mu^{1/2} \hat{S}_{1} + \mu \hat{S}_{2} + \mu^{3/2} \hat{S}_{3})^{2} + \frac{1}{6} \hat{\psi}''''(\hat{S})(\mu^{1/2} \hat{S}_{1} + \mu \hat{S}_{2} + \mu^{3/2} \hat{S}_{3})^{3}
+ \frac{1}{24} \hat{\psi}''''(\hat{S} + \theta(\mu^{1/2} \hat{S}_{1} + \mu \hat{S}_{2} + \mu^{3/2} \hat{S}_{3}))(\mu^{1/2} \hat{S}_{1} + \mu \hat{S}_{2} + \mu^{3/2} \hat{S}_{3})^{4}
= \mu^{1/2} \hat{\psi}''''(\hat{S}) \hat{S}_{1} + \mu(\hat{\psi}''(\hat{S}) \hat{S}_{2} + \frac{1}{2} \hat{\psi}'''(\hat{S}) \hat{S}_{1}^{2})
+ \mu^{3/2}(\hat{\psi}''''(\hat{S}) \hat{S}_{1} \hat{S}_{2} + \frac{1}{6} \hat{\psi}''''(\hat{S}) \hat{S}_{1}^{-}) + \mu^{2} \overline{R}_{\psi}. \]

Here we used that \( \hat{\psi}'(\hat{S}) = 0 \), by (2.38). Equations (5.44) and (5.45) imply that in the domain \( Q \setminus \Gamma \)
\[
(\mu \lambda)^{1/2} \partial_{S} S_{2}^{(\mu)} + c \left( \partial_{S} W(\varepsilon(\nabla_{x} u_{2}^{(\mu)}), S_{2}^{(\mu)}) \right) + \frac{1}{\mu^{1/2}} \hat{\psi}'(S_{2}^{(\mu)}) - \mu^{1/2} \lambda \Delta_{x} S_{2}^{(\mu)} 
= c \left( -\tau : \hat{T} + \hat{\psi}''(\hat{S}) \hat{S}_{1} \right)
+ c \mu^{1/2} \left( -\tau : \hat{T} + \hat{\psi}''(\hat{S}) \hat{S}_{2} + \frac{1}{2} \hat{\psi}'''(\hat{S}) \hat{S}_{1}^{2} \right)
+ c \mu \left( -\tau : \hat{T} + \hat{\psi}''(\hat{S}) \hat{S}_{3} + \hat{\psi}''(\hat{S}) \hat{S}_{1} \hat{S}_{2} + \frac{1}{6} \hat{\psi}''''(\hat{S}) \hat{S}_{1}^{3} - \lambda \Delta_{x} \hat{S}_{1} + \frac{\lambda^{1/2}}{c} \partial_{t} \hat{S}_{1} \right)
+ \mu^{3/2} \overline{R}_{S_{1} + c(\ldots)} , \]
where
\[
\overline{R}_{S_t+c(...)} = c ( \varpi : D \varpi \tilde{S}_3 + \overline{R}_0 - \lambda \Delta_x (\tilde{S}_2 + \mu^{1/2} \tilde{S}_3)) + \lambda^{1/2} \partial_t (\tilde{S}_2 + \mu^{1/2} \tilde{S}_3). \tag{5.47}
\]

Here we used that the function \( \tilde{S} \) has the constant values 0 in \( \gamma \) and 1 in \( \gamma' \).

**Corollary 5.7** Assume that the functions \( \tilde{S}_1, \tilde{S}_2 \) and \( \tilde{S}_3 \) satisfy (3.30) – (3.32). Then there is a constant \( K \) such that for all \( (t, x) \in Q \setminus \Gamma \) and all \( 0 < \mu \leq \mu_0, 0 < \lambda \leq \lambda_0 \)
\[
\left| \partial_t S_2^{(\mu)} + \frac{c}{(\mu \lambda)^{1/2}} \left( \partial_3 W(\varepsilon (\nabla_x u_2^{(\mu)}), S_2^{(\mu)}) + \frac{1}{\mu^{1/2}} \hat{\gamma}'(S_2^{(\mu)}) - \mu^{1/2} \lambda \Delta_x S_2^{(\mu)} \right) \right| \leq \frac{K \mu}{\lambda^{1/2}}. \tag{5.48}
\]

**Proof:** By (3.30) – (3.32), the brackets on the right hand side of equation (5.46) vanish. Therefore, if we divide the latter equation by \( (\mu \lambda)^{1/2} \), we obtain
\[
\partial_t S_2^{(\mu)} + \frac{c}{(\mu \lambda)^{1/2}} \left( \partial_3 W(\varepsilon (\nabla_x u_2^{(\mu)}), S_2^{(\mu)}) + \frac{1}{\mu^{1/2}} \hat{\gamma}'(S_2^{(\mu)}) - \mu^{1/2} \lambda \Delta_x S_2^{(\mu)} \right) = \frac{\mu}{\lambda^{1/2}} \overline{R}_{S_t+c(...)}.
\]
\( \tilde{S}_1 \) and \( \tilde{S}_2 \) are given in (3.36), (3.37), and the function \( \tilde{S}_3 \) is obtained by solving (3.32) for this function. From these equations we see by our general regularity assumptions that \( \| (\tilde{S}_1, \tilde{S}_2, \tilde{S}_3) \|_{L^\infty(Q \setminus \Gamma)} \leq K_1 \), with the constant \( K_1 \) independent of \( \mu \). Using this, we see by inspection of every term in (5.47) that \( \| \overline{R}_{S_t+c(...)} \|_{L^\infty(Q \setminus \Gamma)} \leq K \), with \( K \) independent of \( \mu \) and \( \lambda \). This inequality and the equation above imply (5.48).

### 5.3 Auxiliary estimates needed in the matching region

The following auxiliary estimates are needed to prove (2.51) and (2.53) in the matching region \( Q_{\text{match}}^{(\mu)} \).

**Lemma 5.8** The functions \( u_2^{(\mu)} \), \( T_2^{(\mu)} \), \( S_2^{(\mu)} \) defined in (3.25) – (3.27) and \( u_1^{(\mu)} \), \( T_1^{(\mu)} \), \( S_1^{(\mu)} \) defined in (3.49) – (3.51) satisfy
\[
\| S_1^{(\mu)} - S_2^{(\mu)} \|_{L^\infty(Q_{\text{match}}^{(\mu)})} \leq K \mu^{3/2} |\ln \mu|^2, \tag{5.49}
\]
\[
\| D_x^\alpha (S_1^{(\mu)} - S_2^{(\mu)}) \|_{L^\infty(Q_{\text{match}}^{(\mu)})} \leq K \lambda^{-|\alpha|/2} \mu^{3-|\alpha|/2}, \quad 1 \leq |\alpha| \leq 2, \tag{5.50}
\]
\[
\| \partial_t (S_1^{(\mu)} - S_2^{(\mu)}) \|_{L^\infty(Q_{\text{match}}^{(\mu)})} \leq K \lambda^{-1/2} \mu, \tag{5.51}
\]
\[
\| u_1^{(\mu)} - u_2^{(\mu)} \|_{L^\infty(Q_{\text{match}}^{(\mu)})} \leq K \lambda^{1/2} \mu^{3/2} |\ln \mu|, \tag{5.52}
\]
\[
\| \nabla_x (u_1^{(\mu)} - u_2^{(\mu)}) \|_{L^\infty(Q_{\text{match}}^{(\mu)})} \leq K \mu, \tag{5.53}
\]
\[
\| T_1^{(\mu)} - T_2^{(\mu)} \|_{L^\infty(Q_{\text{match}}^{(\mu)})} \leq K \mu, \tag{5.54}
\]
for all \( \mu \in (0, \mu_0], \lambda \in (0, \lambda_0] \). Here \( \alpha \) denotes a multi-index and \( K \) denotes a positive constant, which does not necessarily have the same value in the six estimates.

**Proof:** Since the proofs of these estimates are long and technical, we present here only the proofs of (5.49) and (5.52). The proofs of the estimates (5.51), (5.53) and (5.54) run
along the same lines. (5.54) is an immediate consequence of the definitions (3.37), (3.51) of $T_2^{(\mu)}$ and $T_1^{(\mu)}$, and of the estimates (5.53), (5.49).

In this proof we mostly drop the arguments $t$ and $\eta$ to simplify the notation. As usual we write $\zeta = \frac{\xi}{(\mu \lambda)^{1/2}}$. We need that for $(t, \eta, \xi) \in Q_{\text{match}}^{(\mu \lambda)}$ the inequalities

$$\frac{3}{2} |\ln \mu| a \leq |\frac{\xi}{(\mu \lambda)^{1/2}}| = |\zeta| \leq \frac{3}{2} |\ln \mu| a, \quad (5.55)$$

hold, by definition of $Q_{\text{match}}^{(\mu \lambda)}$ in (2.34).

We begin with the proof of (5.49). By definition of $S_1^{(\mu)}$ and $S_2^{(\mu)}$ in (3.50) and (3.26) we have

$$|S_1^{(\mu)} - S_2^{(\mu)}| = |S_0 + \mu^{1/2}S_1 + \mu S_2 - \hat{S} - \mu^{1/2} \hat{S}_1 - \mu \hat{S}_2 - \mu^{3/2} \hat{S}_3| \leq |S_0 - \hat{S}| + \mu^{1/2}|S_1 + \mu^{1/2}S_2 - \hat{S}_1 - \mu^{1/2} \hat{S}_2| + \mu^{3/2}|\hat{S}_3|. \quad (5.56)$$

To estimate the first term on the right hand side note that since $\hat{S}(t, x) = \hat{S}(\xi) = 1^+(\xi)$, relations (3.33), (3.31), and (5.55) imply for $(t, \eta, \xi) \in Q_{\text{match}}^{(\mu \lambda)}$ that

$$|S_0 \left(\frac{\xi}{\mu \lambda}^{1/2}\right) - \hat{S}(\xi)| \leq K_1 e^{-a|\xi|} \leq K_1 e^{-\frac{3}{2}|\ln \mu|} = K_1 \mu^{3/2}. \quad (5.57)$$

To estimate the second term on the right hand side of (5.56) we introduce the notations

$$\tilde{\rho}_1(\zeta) = \begin{cases} \frac{\pi \hat{T}(\zeta)}{\psi(0)} & \zeta < 0, \\
\frac{\pi \hat{T}(\zeta)}{\psi(0)} & \zeta > 0,
\end{cases} \quad (5.58)$$

$$\tilde{\rho}_2(\zeta) = \begin{cases} \frac{1}{\psi(0)} \left( \frac{\hat{T}(-)}{2} - \frac{\hat{\phi}''(0)}{2} \left( \frac{\pi \hat{T}(\zeta)}{\psi(0)} \right)^2 + \lambda \frac{\hat{\phi}''(0)}{2} \right) & \zeta < 0, \\
+ \lambda \frac{\hat{T}(\zeta)}{2} & \zeta > 0.
\end{cases} \quad (5.59)$$

By definition of the functions $\rho_1, \rho_2$ in (4.10) and (3.59), and by definition of $\varphi, \varphi_{\pm}$ in (3.57) and (3.58), we have

$$\tilde{\rho}_i(\zeta) - \rho_i(\zeta) = (1 - \varphi(-\zeta) - \varphi(\zeta)) \tilde{\rho}_i(\zeta) = 0, \quad \text{for } i = 1, 2 \text{ and } |\zeta| \geq 2.$$
\(3.36\), \(3.37\), and employ the triangle inequality to obtain

\[
|S_1 + \mu^{1/2}S_2 - \tilde{S}_1 - \mu^{1/2}\tilde{S}_2| \\
\leq \left| \tilde{\rho}_1 + \mu^{1/2}\tilde{\rho}_2 - \frac{\varepsilon}{\psi''(\tilde{S})} + \mu^{1/2} \left( \frac{\varepsilon}{\psi''(\tilde{S})} - \frac{\psi''(\tilde{S})}{2\psi''(\tilde{S})} \left( \frac{\varepsilon}{\psi''(\tilde{S})} \right) \right)^2 \right| \\
+ |S_1 - \tilde{\rho}_1| + \mu^{1/2}|S_2 - \tilde{\rho}_2| = |I_1| + |I_2| + |I_3|. 
\tag{5.62}
\]

By \(5.60\), \(5.61\) and \(5.55\) we have for \((t, \eta, \xi) \in Q_{\text{match}}^{(\mu\lambda)}\)

\[
|I_2| + |I_3| = |S_1(t, \eta, \xi) - \tilde{\rho}_1(t, \eta, \xi)| + \mu^{1/2}|S_2(t, \eta, \xi) - \tilde{\rho}_2(t, \eta, \xi)| \\
\leq (K_2 + \bar{K})e^{-a|\xi|} + \mu^{1/2}(K_3(1 + |\xi|) + \bar{K})e^{-a|\xi|} \\
\leq \left( C_1 + \mu^{1/2}C_2 \left( 1 + 3|\ln\mu| \right) \right) e^{-\frac{a}{2} |\ln\mu|} \leq C_3 \mu^{3/2}. \tag{5.63}
\]

To find an estimate for \(|I_1|\) note that the definitions of \(\tilde{\rho}_1, \tilde{\rho}_2\) in \((5.58), (5.59)\) yield for \((t, \eta, \xi) \in Q_{\text{match}}^{(\mu\lambda)}\) with \(\xi > 0\)

\[
I_1 = \tilde{\rho}_1 \left( \frac{\xi}{(\mu\lambda)^{1/2}} \right) + \mu^{\frac{1}{4}}\tilde{\rho}_2 \left( \frac{\xi}{(\mu\lambda)^{1/2}} \right) - \frac{\varepsilon}{\psi''(1)} - \mu^{\frac{1}{4}} \left( \frac{\varepsilon}{\psi''(1)} - \frac{\psi''(1)}{2\psi''(1)} \left( \frac{\varepsilon}{\psi''(1)} \right)^2 \right) \\
= \frac{1}{\psi''(1)} \left( \varepsilon: \hat{T}^{(+) - \varepsilon: \hat{T}(\xi)} + \sigma'(0)\xi + \varepsilon: De(a^* \otimes n + \nabla_{\Gamma} u^*)\xi \right) \\
+ \mu^{\frac{1}{4}}(\varepsilon: \hat{T}^{(+) - \varepsilon: \hat{T}(\xi)} - \mu^{\frac{1}{4}} \frac{\psi''(1)}{2} \left( \frac{\varepsilon}{\psi''(1)} \right)^2 \left( \frac{\varepsilon}{\psi''(1)} \right) \right). \tag{5.64}
\]

\((5.2)\) and \((5.55)\) together yield

\[
|\varepsilon: \hat{T}^{(+) - \varepsilon: \hat{T}(\xi)} + \sigma'(0)\xi + \varepsilon: De(a^* \otimes n + \nabla_{\Gamma} u^*)\xi - \varepsilon: \hat{T}(\xi)| \\
= |R_{\varepsilon, \hat{T}}(\xi)\xi^2| \leq C_4 \mu\lambda |\ln\mu|^2. \tag{5.65}
\]

Since \(\hat{T}^{(+) = \hat{T}(0+)\), \(\hat{T}^{(+) = \hat{T}(0+)}\), the mean value theorem and \((5.55)\) imply

\[
\mu^{1/2} \left| \varepsilon: \left( \hat{T}^{(+) - \hat{T}(\xi)} \right) - \frac{\psi''(1)}{2(\psi''(1))^2} \left( \varepsilon: \hat{T}^{(+) - (\varepsilon: \hat{T}(\xi))} \right) \right| \\
= \mu^{1/2} |R(\xi)| \leq C_5 \mu\lambda^{1/2} |\ln\mu|, \tag{5.66}
\]

where the remainder term \(R) \in L^\infty(U_\delta).\) Combination of \((5.64) - (5.66)\) yields for \((t, \eta, \xi) \in Q_{\text{match}}^{(\mu\lambda)}\) with \(\xi > 0\) that

\[
|I_1| \leq \mu(C_4 \lambda |\ln\mu|^2 + C_5 \lambda^{1/2} |\ln\mu|) \leq C_6 \lambda^{1/2} |\ln\mu|^2 \mu. \tag{5.67}
\]

From the definitions of \(\tilde{\rho}_1, \tilde{\rho}_2\) in \((5.58), (5.59)\) we see that for \((t, \eta, \xi) \in Q_{\text{match}}^{(\mu\lambda)}\) with \(\xi < 0\) the term \(I_1\) takes the form

\[
I_1 = \frac{1}{\psi''(1)} \left( \varepsilon: \hat{T}^{(-)} - \varepsilon: \hat{T}(\xi) + \sigma'(0)\xi \right) \\
+ \mu^{1/2} \left( \varepsilon: \hat{T}^{(-)} - \varepsilon: \hat{T}(\xi) - \mu^{1/2} \frac{\psi''(0)}{2} \left( \frac{\varepsilon}{\psi''(0)} \right)^2 \left( \frac{\varepsilon}{\psi''(0)} \right) \right). 
\]

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Using \((5.3)\) instead of \((5.2)\), we see as above that the estimate \((5.67)\) also holds in this case, whence the estimate \((5.67)\) is valid for all \((t, \eta, \xi) \in Q_{\text{match}}^{(\mu\lambda)}\).

To finish the proof of \((5.49)\), we combine \((5.56)\) with \((5.57), (5.62), (5.63)\) and \((5.67)\) to obtain the estimate

\[
|S_2^{(\mu)} - S_2^{(\mu)}| \leq K_1 \mu^{3/2} + \mu^{1/2} (C_6 \lambda^{1/2} |\ln \mu|^2 \mu + C_3 \mu^{3/2}) + \mu^{3/2} |\tilde{S}_3|,
\]

which implies \((5.49)\).

Next we prove \((5.52)\). From \((3.25)\) and \((3.38), (3.39)\) we conclude for \((t, x) \in U_0\) that

\[
u_2^{(\mu)}(t, x) = \tilde{u}(t, x) + \mu^{1/2} \tilde{u}(t, x) + \mu \tilde{u}(t, x)
= u^* \xi + a^* \frac{1}{2} (\xi^*)^2 + \mu^{1/2} u^* \left( \frac{\xi}{\psi''(1)} + \frac{\xi}{\psi''(0)} \right)
+ \tilde{\varphi}(t, x) + \mu^{1/2} \tilde{\varphi}(t, x) + \mu \tilde{u}(t, x).
\]

Combination of this equation with \((3.49)\) and insertion of \((5.52) - (5.54)\) yields with \(\zeta = \frac{\xi}{(\mu \lambda)^{1/2}}\) that

\[
u_1^{(\mu)}(t, x) - \nu_2^{(\mu)}(t, x)
= (\mu \lambda)^{1/2} u_0 \left( \frac{\xi}{(\mu \lambda)^{1/2}} \right) + \mu \lambda^{1/2} u_1 \left( \frac{\xi}{(\mu \lambda)^{1/2}} \right) + \mu \lambda u_2 \left( \frac{\xi}{(\mu \lambda)^{1/2}} \right)
- \left( u^* \xi + a^* \frac{1}{2} (\xi^*)^2 + \mu^{1/2} u^* \left( \frac{\xi}{\psi''(1)} + \frac{\xi}{\psi''(0)} \right) + \mu \tilde{u}(\xi) \right)
= (\mu \lambda)^{1/2} u^* \left( \int_{-\infty}^{\zeta} S_0(\vartheta) d\vartheta - \zeta^* \right) + \mu \lambda a^* \int_{-\infty}^{\zeta} \left( \int_{-\infty}^{\vartheta} S_0(\vartheta_1) d\vartheta_1 - \vartheta^* \right) d\vartheta
+ \mu \lambda^{1/2} u^* \int_{0}^{\zeta} \left( S_1(\vartheta) - \frac{\xi}{\psi''(1)} 1^+(\vartheta) - \frac{\xi}{\psi''(0)} 1^- (\vartheta) \right) d\vartheta - \mu \tilde{u}(\xi)
= (\mu \lambda)^{1/2} J_1(\zeta) + \mu \left( J_2(\zeta) + J_3(\zeta) - \tilde{u}(\xi) \right).
\]

To estimate the right hand side we use the boundary condition \((3.35)\) for \(\tilde{u}\), note that by \((5.58)\) the equation \(\tilde{p}_1(\zeta) = \frac{\xi}{\psi''(1)}\) holds for \(\zeta > 0\), and employ the inequalities \((5.60), (4.8)\) to compute for \((t, \eta, \xi) \in Q_{\text{match}}^{(\mu\lambda)}\) with \(\zeta > 0\)

\[
|J_2(\zeta) + J_3(\zeta) - \tilde{u}(\vartheta)| = |J_2(\zeta) - \lambda a^* \int_{-\infty}^{\vartheta} S_0(\vartheta_1) d\vartheta_1 - \vartheta^*| d\vartheta
+ J_3(\zeta) - \lambda^{1/2} u^* \int_{0}^{\zeta} \left( S_1(\vartheta) - \frac{\xi}{\psi''(1)} 1^+(\vartheta) - \frac{\xi}{\psi''(0)} 1^- (\vartheta) \right) d\vartheta
= |\lambda a^* \int_{-\infty}^{\vartheta} S_0(\vartheta_1) d\vartheta_1 - \vartheta^*| d\vartheta
+ \lambda^{1/2} u^* \int_{0}^{\zeta} \left( S_1(\vartheta) - \tilde{p}_1(\zeta) \right) d\vartheta
\leq |\lambda a^*| \int_{-\infty}^{\vartheta} \frac{K_1}{\vartheta} e^{-\vartheta} d\vartheta + \lambda^{1/2} u^* \int_{0}^{\zeta} (K_2 + \tilde{K}) e^{-\vartheta} d\vartheta
= \left( |\lambda a^*| \frac{K_1}{\vartheta} + \lambda^{1/2} u^* \frac{K_2 + \tilde{K}}{\vartheta} \right) e^{-\vartheta} \leq C_1 \lambda^{1/2} (\|a^*\|_{L^\infty(\Gamma)} + \|u^*\|_{L^\infty(\Gamma)}) e^{-\vartheta}.
\]

(5.69)
The mean value theorem implies for \( \xi > 0 \)
\[
\tilde{u}(t, x) = \tilde{u}(t, \eta, \xi) = \tilde{u}(t, \eta, 0+) + R_\tilde{u}(t, \eta, \xi). 
\]
Since \( \tilde{u}(t, \eta, 0+) = \tilde{u}^{(+)}(t, \eta) \), we infer from this equation and from (5.69), (5.65) for all \( (t, \eta, \xi) \in Q_{\text{match}}^{(\mu \lambda)} \) with \( \xi > 0 \) that
\[
|J_2(\xi) + J_3(\xi) - \tilde{u}(\xi)| \leq |J_2(\xi) + J_3(\xi) - \tilde{u}^{(+)}| + |\tilde{u}(\xi) - \tilde{u}^{(+)}| 
\]
\[
\leq C_1 \lambda^{1/2} \langle \| a^* \| \rangle_{L^\infty(\Gamma)} + \| u^* \|_{L^\infty(\Gamma)} e^{-a \xi} + \| R_{\tilde{u}} \|_{L^\infty(\Omega_0)} \xi 
\]
\[
\leq C_2 \lambda^{1/2} \langle \| a^* \| \rangle_{L^\infty(\Gamma)} + \| u^* \|_{L^\infty(\Gamma)} e^{-3 \frac{1}{2} \ln \mu} + \| R_{\tilde{u}} \|_{L^\infty(\Omega_0)} \xi \frac{3}{a} (\mu \lambda)^{1/2} | \ln \mu | 
\]
\[
\leq C_2 \lambda^{1/2} \mu^{1/2} | \ln \mu |. 
\]

Using the boundary condition (3.34) for \( \tilde{u} \) instead of (3.35), we see by the analogous computation that (5.70) also holds for \( (t, \eta, \xi) \in Q_{\text{match}}^{(\mu \lambda)} \) with \( \xi < 0 \).

Now use (5.70) to estimate the second term on the right hand side of (5.68). The first term is estimated by (1.8). Because (5.55) yields \( e^{-a |\xi|} \leq e^{-\frac{3}{2} \ln \mu} = \mu^{3/2} \), we obtain for \( (t, \eta, \xi) \in Q_{\text{match}}^{(\mu \lambda)} \)
\[
|u_1^{(\mu)}(t, x) - u_2^{(\mu)}(t, x)| \leq (\mu \lambda)^{1/2} |u^*| \int_{-\infty}^\xi S_0(\theta) d\theta + \mu |J_2(\xi) + J_3(\xi) - \tilde{u}(\xi)| 
\]
\[
\leq (\mu \lambda)^{1/2} \left( \max_{\Gamma} |u^*| \right) \frac{K_1}{a} e^{-a |\xi|} + C_2 \lambda^{1/2} \mu^{3/2} | \ln \mu | \leq C_1 \lambda^{1/2} \mu^2 + C_2 \lambda^{1/2} \mu^{3/2} | \ln \mu |, 
\]
which implies (5.52).

5.4 End of the proof of Theorem 2.3

To complete the proof of Theorem 2.3, note first that (3.22), (3.23) imply
\[
u^{(\mu)}|_{\partial \Omega} = u_2^{(\mu)}|_{\partial \Omega}, \quad \partial_{\nu_{\partial \Omega}} S_1^{(\mu)}|_{\partial \Omega} = \partial_{\nu_{\partial \Omega}} S_2^{(\mu)}|_{\partial \Omega}, \quad S_1^{(\mu)}|_{\Omega_{\text{inn}}} = S_1^{(\mu)}|_{\Omega_{\text{match}}}.
\]

Therefore (2.50) follows from the definition of \( S_1^{(\mu)} \) in (3.51), equation (2.45) follows from (2.28), (2.31), (3.60), and (2.49) is a consequence of the definition of \( u_2^{(\mu)}|_{\partial \Omega} \) in (3.25) and of (2.20), (2.25), (3.33). Moreover, the estimate (2.55) for the right hand side \( f_3^{(\mu \lambda)} \) of (2.50) follows from the definition of \( S_2^{(\mu)}|_{\partial \Omega} \) in (3.26) and from \( \partial_{n_{\partial \Omega}} \hat{S}|_{\partial \Omega} = 0 \). This last equation holds, since by assumption \( \Gamma(t) \subseteq \Omega \), which implies that \( \hat{S}(t) \) is identically equal to 0 or 1 in a neighborhood of \( \partial \Omega \).

(2.47) follows from the definition of \( T^{(\mu)} \) in (3.24); equation (2.57) is an immediate consequence of (1.1).

It remains to verify the estimates (2.51) – (2.54) for the right hand sides \( f_1^{(\mu \lambda)} \), \( f_2^{(\mu \lambda)} \) of the equations (2.46) and (2.48). To this end we put together all the estimates derived in Sections 5.1 – 5.3. We start with the proof of (2.51) and (2.52).

Equation (3.22) yields
\[
\nabla_x u^{(\mu)} = \nabla_x u_1^{(\mu)} \phi_{\mu \lambda} + \nabla_x u_2^{(\mu)} (1 - \phi_{\mu \lambda}) + \left( u_1^{(\mu)} - u_2^{(\mu)} \right) \otimes \nabla_x \phi_{\mu \lambda}.
\]
We insert this equation into (3.21) and use (3.23) and (3.27), (3.51) to obtain
\[ T^{(\mu)} = T_1^{(\mu)} \phi_{\mu \lambda} + T_2^{(\mu)} (1 - \phi_{\mu \lambda}) + D \varepsilon ((u_1^{(\mu)} - u_2^{(\mu)}) \otimes \nabla \phi_{\mu \lambda}). \]  
(5.71)

The function \( \phi_{\mu \lambda} \) defined in (3.21) is independent of \((t, \eta)\). The decomposition (3.5) of the gradient thus yields
\[ \nabla \phi_{\mu \lambda} = \frac{2a}{3(\lambda \mu)^{1/2} |\ln \mu|} \phi'_{\mu \lambda} n, \]  
(5.72)

with the unit normal vector \( n = n(t, \eta) \) to \( \Gamma_{\xi}(t) \) at \( \eta \in \Gamma_{\xi}(t) \). We write
\[ \phi'_{\mu \lambda} = \phi'(\frac{2a}{3(\mu \lambda)^{1/2} |\ln \mu|}), \quad \phi''_{\mu \lambda} = \phi''(\frac{2a \xi}{3(\mu \lambda)^{1/2} |\ln \mu|}), \]  
by a slight abuse of notation. With (5.71) and (5.72) we compute
\[ \text{div}_x T^{(\mu)} + b = (\text{div}_x T_1^{(\mu)} + b) \phi_{\mu \lambda} + (\text{div}_x T_2^{(\mu)} + b)(1 - \phi_{\mu \lambda}) \]
\[ + \left( (T_1^{(\mu)} - T_2^{(\mu)}) n + \text{div}_x D \varepsilon ((u_1^{(\mu)} - u_2^{(\mu)}) \otimes n) \right) \frac{2a}{3(\lambda \mu)^{1/2} |\ln \mu|} \phi'_{\mu \lambda} n \]
\[ + \left( D \varepsilon ((u_1^{(\mu)} - u_2^{(\mu)}) \otimes n) \right) n \left( \frac{2a}{3(\lambda \mu)^{1/2} |\ln \mu|} \right)^2 \phi''_{\mu \lambda}. \]  
(5.73)

Inequality (2.52) is an immediate consequence of this equation and of (5.29), since \( \phi_{\mu \lambda} = 0 \) in \( Q_{\text{out}}^{(\mu \lambda)} \), and (2.51) is obtained by estimating the right hand side of (5.73) using the obvious inequality
\[ |\text{div}_x D \varepsilon ((u_1^{(\mu)} - u_2^{(\mu)}) \otimes n)| \leq C (|\nabla ((u_1^{(\mu)} - u_2^{(\mu)})| + |u_1^{(\mu)} - u_2^{(\mu)})|) \]  
(5.74)

and the equation and inequalities (5.21), (5.29), (5.52) – (5.54).

We next proof (2.53) and (2.54). The inequality (2.54) follows immediately from (5.48), since \( \phi_{\mu \lambda} = 0 \) on \( Q_{\text{out}}^{(\mu \lambda)} \), which by (3.22) and (3.23) implies \( (u^{(\mu)}, S^{(\mu)}) = (u_2^{(\mu)}, S_2^{(\mu)}) \) on \( Q_{\text{out}}^{(\mu \lambda)} \subseteq Q \setminus \Gamma \).

It remains to verify (2.53). Since \( W_S(\varepsilon, S) = -\varepsilon : D(\varepsilon - \varepsilon S) \), by (1.8), it follows from (3.24), (3.51), and (5.71) that
\[ W_S(\varepsilon (\nabla_x u^{(\mu)}), S^{(\mu)}) - W_S(\varepsilon (\nabla_x u_1^{(\mu)}), S_1^{(\mu)}) = -\varepsilon : (T^{(\mu)} - T_1^{(\mu)}) \]
\[ = -\varepsilon : (T_2^{(\mu)} - T_1^{(\mu)})(1 - \phi_{\mu \lambda}) - \varepsilon : D \varepsilon ((u_1^{(\mu)} - u_2^{(\mu)}) \otimes \nabla \phi_{\mu \lambda}). \]  

The mean value theorem and (3.23) imply
\[ \hat{\psi}'(S^{(\mu)}) - \hat{\psi}'(S_1^{(\mu)}) = \hat{\psi}'(S_1^{(\mu)} + \vartheta(S_2^{(\mu)} - S_1^{(\mu)})(1 - \phi_{\mu \lambda}))(S_2^{(\mu)} - S_1^{(\mu)})(1 - \phi_{\mu \lambda}), \]  
for a suitable \( 0 < \vartheta(t, x) < 1 \), and
\[ \Delta_x S^{(\mu)} - \Delta_x S_1^{(\mu)} = \Delta_x (S_2^{(\mu)} - S_1^{(\mu)})(1 - \phi_{\mu \lambda}) + 2 \nabla_x (S_1^{(\mu)} - S_2^{(\mu)}) \cdot \nabla \phi_{\mu \lambda} \]
\[ + (S_1^{(\mu)} - S_2^{(\mu)}) \Delta_x \phi_{\mu \lambda}. \]  
(5.75)
The right hand sides of the last three equations vanish on the set $Q_{\text{inn}}^{(\mu, \lambda)}$, since $\phi_{\mu \lambda} = 1$ on this set. On the set $Q_{\text{match}}^{(\mu, \lambda)}$, the right hand sides of these equations can be estimated using (5.49), (5.50), (5.52), (5.54), (5.72). In the estimation of (5.75) we also note that since $\phi_{\mu \lambda}$ is independent of $(t, \eta)$, analogous to (5.35) the equation

$$
\Delta_x \phi_{\mu \lambda} = -\kappa \partial_t \phi_{\mu \lambda} + \partial_x^2 \phi_{\mu \lambda} = -\kappa \frac{2a}{3(\lambda \mu)^{1/2} |\ln \mu|} \phi'_{\mu \lambda} + \left( \frac{2a}{3(\lambda \mu)^{1/2} |\ln \mu|} \right)^2 \phi''_{\mu \lambda}
$$

holds, with twice the mean curvature $\kappa(t, \eta, \xi)$ of the surface $\Gamma_\xi(t)$ at $\eta \in \Gamma_\xi(t)$. Together we obtain that on $Q_{\text{inn}}^{(\mu, \lambda)} \cup Q_{\text{match}}^{(\mu, \lambda)}$ the inequality

$$
\left| W_S(\varepsilon(\nabla_x u^{(\mu)}), S^{(\mu)}) + \frac{1}{\mu^{1/2}} \psi'(S^{(\mu)}) - \mu^{1/2} \lambda \Delta_x S^{(\mu)} \right|
$$

$$
\left| W_S(\varepsilon(\nabla_x u_1^{(\mu)}), S_1^{(\mu)}) + \frac{1}{\mu^{1/2}} \psi'(S_1^{(\mu)}) - \mu^{1/2} \lambda \Delta_x S_1^{(\mu)} \right|
$$

$$
\leq K(\mu + \mu |\ln \mu|^2 + \mu^{1/2} \lambda \mu^{1/2} \lambda^{-1}) \leq K|\ln \mu|^2 \quad \text{(5.76)}
$$

holds. Similarly, (5.23) implies

$$
\partial_t S^{(\mu)} - \partial_t S_1^{(\mu)} = \partial_t (S_2^{(\mu)} - S_1^{(\mu)}) (1 - \phi_{\mu \lambda}) + (S_1^{(\mu)} - S_2^{(\mu)}) \partial_t \phi_{\mu \lambda}.
$$

The right hand side of this equation vanishes on $Q_{\text{inn}}^{(\mu, \lambda)}$. To estimate the right hand side on the set $Q_{\text{match}}^{(\mu, \lambda)}$ we use the inequalities (5.49), (5.51) and the equation

$$
\partial_t \phi_{\mu \lambda} = -\frac{2a S^{(\mu)}}{3(\mu \lambda)^{1/2} |\ln \mu|} \phi'_{\mu \lambda},
$$

which follows from (3.21) and Lemma 5.5. The result is

$$
|\partial_t S^{(\mu)} - \partial_t S_1^{(\mu)}| \leq K \lambda^{-1/2} |\ln \mu|, \quad \text{on } Q_{\text{inn}}^{(\mu, \lambda)} \cup Q_{\text{match}}^{(\mu, \lambda)}.
$$

By combination of (5.41), (5.76) and (5.77) we see that the inequality

$$
\left| \partial_t S^{(\mu)} + \frac{c}{(\mu \lambda)^{1/2}} \left( W_S(\varepsilon(\nabla_x u^{(\mu)}), S^{(\mu)}) + \frac{1}{\mu^{1/2}} \psi'(S^{(\mu)}) - \mu^{1/2} \lambda \Delta_x S^{(\mu)} \right) \right|
$$

$$
\left| \partial_t S_1^{(\mu)} + \frac{c}{(\mu \lambda)^{1/2}} \left( W_S(\varepsilon(\nabla_x u_1^{(\mu)}), S_1^{(\mu)}) + \frac{1}{\mu^{1/2}} \psi'(S_1^{(\mu)}) - \mu^{1/2} \lambda \Delta_x S_1^{(\mu)} \right) \right|
$$

$$
\leq K \lambda^{-1/2} |\ln \mu| + \frac{c}{(\mu \lambda)^{1/2}} K \lambda |\ln \mu|^2
$$

$$
\leq K \left( \frac{\mu^2}{\lambda} \right)^{1/2} |\ln \mu|^2 + K \frac{\mu^{1/2}}{\lambda} |\ln \mu| + cK \left( \frac{\mu^2}{\lambda} \right)^{1/2} |\ln \mu|^2 \leq K_1 \left( \frac{\mu^2}{\lambda} \right)^{1/2} |\ln \mu|^2
$$

holds on the set $Q_{\text{inn}}^{(\mu, \lambda)} \cup Q_{\text{match}}^{(\mu, \lambda)}$. This proves (2.53) and completes the proof of Theorem 2.3. \hfill \blacksquare
6 Proof of Theorem 2.8

This section contains the proof of Theorem 2.8. \((u, T, S)\) denotes the asymptotic solution constructed in Theorem 2.3 and \((u_{AC}, T_{AC}, S_{AC})\) denotes the exact solution of (1.1) – (1.3), (2.67) – (2.69). For the proof we need a lemma and a theorem, which we state first.

Lemma 6.1 For all \(x \in \Gamma(\hat{t})\) the propagation speeds \(s_{AC}\) and \(s\) satisfy

\[
s_{AC}(\hat{t}, x) - s(\hat{t}, x) = \frac{1}{|\nabla_x S(\hat{t}, x)|} \left( f_2^{(\mu\lambda)}(\hat{t}, x) - \frac{c}{(\mu\lambda)^{1/2}} \varepsilon : (T_{AC}(\hat{t}, x) - T(\hat{t}, x)) \right),
\]

where \(f_2^{(\mu\lambda)}\) is the right hand side of equation (2.48).

Proof: Since the manifold \(\Gamma\) is a level set of \(S\) and since by (2.70) the manifold \(\Gamma_{AC}\) is a level set of \(S_{AC}\), it follows that \((\partial_t S(\hat{t}, x), \nabla_x S(\hat{t}, x))\) and \((\partial_t S_{AC}(\hat{t}, x), \nabla_x S_{AC}(\hat{t}, x))\) are normal vectors to the respective manifolds at \((\hat{t}, x)\). Moreover, (2.69) implies that \(\nabla_x S_{AC}(\hat{t}, x) = \nabla_x S(\hat{t}, x)\). From (2.9) we thus infer that

\[
s(\hat{t}, x) = -\frac{\partial_t S(\hat{t}, x)}{\nabla_x S(\hat{t}, x) \cdot n(\hat{t}, x)}, \quad s_{AC}(\hat{t}, x) = -\frac{\partial_t S_{AC}(\hat{t}, x)}{\nabla_x S_{AC}(\hat{t}, x)}.\]

(6.2)

(6.3)

For brevity we do not write the argument \((\hat{t}, x)\) in the following computation. In (6.2) we eliminate \(\partial_t S\) with the help of (2.38), and in (6.3) we replace \(\partial_t S_{AC}\) by the right hand side of (1.3). Together with another application of (2.69) this results in

\[
s_{AC} - s = \frac{c}{(\mu\lambda)^{1/2}|\nabla_x S|} \left( \langle \partial_S W(\varepsilon(\nabla_x u_{AC}), S_{AC}) + \frac{1}{\mu^{1/2}} \hat{\psi}'(S) - \mu^{1/2} \lambda \Delta_x S \rangle + \frac{1}{|\nabla_x S|} f_2^{(\mu\lambda)} - \frac{c}{(\mu\lambda)^{1/2}|\nabla_x S|} (\varepsilon : T_{AC} - \varepsilon : T) \right).
\]

which is (6.1). In the last step we used that by (1.8) and (2.47) we have

\[
\partial_S W(\varepsilon(\nabla_x u_{AC}), S_{AC}) = -\varepsilon : T_{AC}, \quad \partial_S W(\varepsilon(\nabla_x u), S) = -\varepsilon : T.
\]

The proof of Lemma 6.1 is complete. \(\blacksquare\)

Theorem 6.2 Suppose that the order of differentiability of \(\hat{\psi}, \Gamma, \hat{u}, \hat{b}\) is higher by two than required in Theorem 2.3. Assume that the principal curvatures \(\kappa_1^{(\mu\lambda)}, \kappa_2^{(\mu\lambda)}\) of the regular \(C^1\) manifold \(\Gamma(\hat{t}) = \Gamma^{(\mu\lambda)}(\hat{t})\) are bounded, uniformly with respect to \(\mu \in (0, \mu_0]\) and \(\lambda \in (0, \lambda_0]\), and that there is an open subset \(\Omega' \subset \subset \Omega\) and \(\delta > 0\) such that the neighborhood \(U_{\delta^{(\mu\lambda)}}^{(\mu\lambda)}(\hat{t})\) of \(\Gamma^{(\mu\lambda)}(\hat{t})\) defined in (2.33) satisfies \(U_{\delta^{(\mu\lambda)}}^{(\mu\lambda)}(\hat{t}) \subseteq \Omega'\). Then there is a constant \(K_5\) such that for all \(\mu \in (0, \mu_0]\) and all \(\lambda \in (0, \lambda_0]\)

\[
\|T_{AC} - T\|_{L^2(\Gamma(\hat{t}))} \leq K_5 |\ln \mu|^3 \mu.
\]

(6.4)
We postpone the proof of this theorem and first finish the proof of Theorem 2.8.

**End of the proof of Theorem 2.8** By (2.7) and (2.56) we have for \( x = (\eta, 0) \in \Gamma(\hat{t}) \) that

\[
\nabla_x S(\hat{t}, x) = n(\partial_n S)(\hat{t}, x) + \nabla_\Gamma S(\hat{t}, x)
\]

\[
= \frac{1}{(\mu \lambda)^{1/2}} \left( S_0'(\zeta) + \mu^{1/2} \partial_\zeta S_1(t, \eta, \zeta) + \mu \partial_\zeta S_2(t, \eta, \zeta) \right) |_{\zeta = 0} n(\hat{t}, \eta) \\
+ \mu^{1/2} \nabla_\Gamma S_1(\hat{t}, \eta, 0) + \mu \nabla_\Gamma S_2(\hat{t}, \eta, 0).
\]

(6.5) implies

\[
S_0'(0) = \sqrt{2\psi(S_0(0))} = \sqrt{2\psi(1/2)} > 0,
\]

whence, from (6.5) for \( \mu \in (0, \mu_0] \) and \( \lambda \in (0, \lambda_0] \) with \( \mu_0 \) sufficiently small,

\[
|\nabla_x S(\hat{t}, x)| \geq \frac{1}{(\mu \lambda)^{1/2}} \left( \sqrt{2\psi(1/2)} - \mu^{1/2} |\partial_\zeta S_1(t, \eta, 0)| - \mu |\partial_\zeta S_2(t, \eta, 0)| \right) \\
- \mu^{1/2} |\nabla_\Gamma S_1(\hat{t}, \eta, 0)| - \mu |\nabla_\Gamma S_2(\hat{t}, \eta, 0)| \geq \frac{1}{2(\mu \lambda)^{1/2}} \sqrt{2\psi(1/2)}. \tag{6.6}
\]

Combination of (6.1) with the inequalities (2.53), (6.4) and (6.6) yields

\[
\|s_{AC}(\hat{t}) - s(\hat{t})\|_{L^2(\Gamma(\hat{t}))} \leq \frac{1}{\min_{\Gamma(\hat{t})} |\nabla_x S(\hat{t})|} \left( \|f_2^{(\mu \lambda)}(\hat{t})\|_{L^2(\Gamma(\hat{t}))} + \frac{c_\Gamma}{(\mu \lambda)^{1/2}} \|T_{AC}(\hat{t}) - T(\hat{t})\|_{L^2(\Gamma(\hat{t}))} \right) \\
\leq \frac{2(\mu \lambda)^{1/2}}{\sqrt{2\psi(1/2)}} \left( |\ln \mu|^2 (\mu \lambda)^{1/2} K_3 \text{meas}(\Gamma(\hat{t}))^{1/2} + \frac{c_\Gamma}{(\mu \lambda)^{1/2}} K_5 |\ln \mu|^3 \mu \right) \\
\leq K_6 |\ln \mu|^2 \mu + K_7 |\ln \mu|^3 \mu.
\]

(2.75) follows from this estimate. The proof of Theorem 2.8 is complete.

**Proof of Theorem 6.2** Note that the function \( (u_{AC}(\hat{t}), T_{AC}(\hat{t})) \) solves the equations (1.1), (1.2) in \( \Omega \) with \( S_{AC}(\hat{t}) = S(\hat{t}) \), by the initial condition (2.69). Moreover, (2.67) holds for \( u_{AC}(\hat{t}) \). From the equations (2.46), (2.47), (2.49) we thus conclude that the difference \( (u_{AC} - u, T_{AC} - T) \) satisfies

\[
-\div_x (T_{AC} - T)(\hat{t}) = -f_1^{(\mu \lambda)}(\hat{t}), \tag{6.7}
\]

\[
(T_{AC} - T)(\hat{t}) = \varepsilon (\nabla_x (u_{AC} - u)(\hat{t})), \tag{6.8}
\]

\[
(u_{AC} - u)(\hat{t}, x) = 0, \quad x \in \partial \Omega. \tag{6.9}
\]

This is the Dirichlet boundary value problem for the elliptic system of elasticity theory in the domain \( \Omega \). It suggests itself to derive the inequality (6.4) by using the \( L^2 \)–regularity theory of elliptic systems, which allows to estimate the norm \( \|T_{AC} - T\|_{L^2(\Gamma(\hat{t}))} \) by the \( L^2 \)–norm of the right hand side \(-f_1^{(\mu \lambda)}(\hat{t})\) of (6.7). To apply this theory directly we would need that the \( L^2 \)–norm of \( f_1^{(\mu \lambda)}(\hat{t}) \) decays to zero for \( \mu \to 0 \) uniformly with respect to \( \lambda \).
However, the relation $\text{meas}(Q_{\text{inn}}^{(\mu\lambda)}(\hat{t}) \cup Q_{\text{match}}^{(\mu\lambda)}(\hat{t})) \leq C_1(\mu\lambda)^{1/2} |\ln \mu|$, which follows from (2.34), and the estimates (2.51), (2.52) yield

$$
\|f_1^{(\mu\lambda)}(\hat{t})\|_{L^2(\Omega)} \leq \left| \ln \mu \right|^{5/2} \frac{\mu^{3/4}}{\lambda^{1/4}} K_1 C_1^{1/2}.
$$

The right hand side does not decay to zero for $\mu \to \infty$ uniformly with respect to $\lambda$, but blows up for $\lambda \to 0$. Therefore direct application of the $L^2$–regularity theory is not possible. Before giving the detailed proof of (6.4) we sketch how to circumvent this difficulty.

Set

$$
A(\mu) = \frac{3}{a} |\ln \mu|^{1/2},
$$

where $a > 0$ is the constant defined in (2.38). By (2.34) we have

$$
Q_{\text{inn}}^{(\mu\lambda)}(\hat{t}) \cup Q_{\text{match}}^{(\mu\lambda)}(\hat{t}) = \left\{ (\eta, \xi) \in U_{\delta}(\hat{t}) \mid |\xi| \leq A(\mu) \lambda^{1/2} \right\}.
$$

Define

\[
\begin{align*}
    f_{11}^{(\mu\lambda)}(x) &= \begin{cases} 
        f_{1}^{(\mu\lambda)}(\hat{t}, x), & x \in Q_{\text{inn}}^{(\mu\lambda)}(\hat{t}) \cup Q_{\text{match}}^{(\mu\lambda)}(\hat{t}), \\
        0, & x \in Q_{\text{out}}^{(\mu\lambda)}(\hat{t}),
    \end{cases} \\
    f_{12}^{(\mu\lambda)}(x) &= \begin{cases} 
        0, & x \in Q_{\text{inn}}^{(\mu\lambda)}(\hat{t}) \cup Q_{\text{match}}^{(\mu\lambda)}(\hat{t}), \\
        f_{1}^{(\mu\lambda)}(\hat{t}, x), & x \in Q_{\text{out}}^{(\mu\lambda)}(\hat{t}),
    \end{cases}
\end{align*}
\]

hence $f_1^{(\mu\lambda)}(\hat{t}) = f_{11}^{(\mu\lambda)} + f_{12}^{(\mu\lambda)}$. For $x = x(\hat{t}, \eta, \xi) \in U_{\delta}(\hat{t})$ we write as usual $f_{11}^{(\mu\lambda)}(x) = f_{11}^{(\mu\lambda)}(\eta, \xi)$. From (2.51) and (6.11) we obtain for $\eta \in \Gamma(\hat{t})$ and $-A(\mu) \lambda^{1/2} \leq \xi \leq A(\mu) \lambda^{1/2}$ that

$$
|f_{11}^{(\mu\lambda)}(\eta, \xi)| \leq |\ln \mu|^2 \left( \frac{\mu}{\lambda} \right)^{1/2} K_1.
$$

Define $\delta_x^{(\mu\lambda)} : \Gamma(\hat{t}) \to \mathbb{R}$ by

$$
\delta_x^{(\mu\lambda)}(\eta) = \int_{-A(\mu)}^{A(\mu)} \lambda^{1/2} f_{11}^{(\mu\lambda)}(\eta, \lambda^{1/2} \xi) \, d\xi.
$$

(6.13) and (6.10) together imply that

$$
|\delta_x^{(\mu\lambda)}(\eta)| \leq 2A(\mu) |\ln \mu|^2 \left( \frac{\mu}{\lambda} \right)^{1/2} K_1 = \frac{6}{a} K_1 |\ln \mu|^3,
$$

for all $\eta \in \Gamma(\hat{t})$. Examination of the boundary value problem (6.7) – (6.9) suggests that for $\mu$ fixed and $\lambda \to 0$ the solution $(u_{AC} - u, T_{AC} - T)$ converges to the solution $(u_*, T_*) : \Omega \to \mathbb{R}^3 \times S^3$ of the transmission problem

\[
\begin{align*}
    -\text{div}_x T_* &= 0, \quad (6.16) \\
    T_* &= D\varepsilon(\nabla_x u_*), \quad (6.17) \\
    [T_*]n &= \delta_x^{(\mu)}, \quad \text{on } \Gamma(\hat{t}), \quad (6.18) \\
    [u_*] &= 0, \quad \text{on } \Gamma(\hat{t}), \quad (6.19) \\
    u_*(x) &= 0, \quad x \in \partial\Omega, \quad (6.20)
\end{align*}
\]
where $\delta^\mu(\eta) = \lim_{\lambda \to 0} \delta^{(\mu\lambda)}(\eta)$ for $\eta \in \Gamma(\hat{t})$. If this limit exists, it follows from (6.15) that

$$|\delta^\mu(\eta)| \leq \frac{6}{a}K_1|\ln \mu|^3.$$ 

This implies that the solution $(u_*, T_*)$ will be bounded by $C\mu|\ln \mu|^3$ with a suitable constant $C$, and this limit behavior suggests that though the $L^2$--norm of $f^\mu(\hat{t})$ blows up for $\lambda \to 0$, the solution $(u_{AC} - u, T_{AC} - T)(\hat{t})$ of (6.7) - (6.9) is bounded by $C|\ln \mu|^3$ with $C$ independent of $\lambda$. The reason for the blow up of $\|f^\mu(\hat{t})\|_{L^2(\Omega)}$ for $\lambda \to 0$ is therefore not that the norm of the solution $(u_{AC} - u, T_{AC} - T)(\hat{t})$ would blow up, but that the solution loses regularity in a neighborhood of the surface $\Gamma(\hat{t})$, which is shown by the equation (6.18) for the limit solution. This equation implies that $T_*$ does not belong to the Sobolev space $W^{1,2}(\Omega)$.

In the following proof we do not study the limit $(u_*, T_*)$. Instead, based on the idea of the behavior of the regularity of $(u_{AC} - u, T_{AC} - T)(\hat{t})$, we decompose this function in the form

$$(u_{AC} - u, T_{AC} - T)(\hat{t}) = (u^\lambda, T^\lambda) + (u^*, T^*)$$

where $(u^\lambda, T^\lambda) = (u^\mu, T^\mu)$ is bounded by $C|\ln \mu|^3$, uniformly with respect to $\lambda$, and for $\lambda \to 0$ has the same regularity behavior as $(u_{AC} - u, T_{AC} - T)$, but otherwise does not approximate $(u_{AC} - u, T_{AC} - T)$. The construction is such that the difference $(u^*, T^*) = (u^\mu_*, T^\mu_*) = (u_{AC} - u, T_{AC} - T) - (u^\lambda, T^\lambda)$ does not lose its regularity for $\lambda \to 0$. Hence, we can use the standard $L^2$--theory for elliptic equations to show that also $(u^*, T^*)$ is bounded by $C|\ln \mu|^3$ independently of $\lambda$.

To construct $(u^\lambda, T^\lambda)$ let $U_\delta(\hat{t})$ be the neighborhood of $\Gamma(\hat{t})$ defined in (2.3) and let $\phi_\delta \in C^\infty_0((-\delta, \delta))$ be a function satisfying

$$\phi_\delta(\xi) = 1, \quad -\delta/2 \leq \xi \leq \delta/2.$$  

We set

$$u^\lambda(x) = \begin{cases} \lambda^{1/2}V(\lambda, \eta, \frac{\xi}{\lambda^{1/2}})\phi_\delta(\xi), & x = x(\hat{t}, \eta, \xi) \in U_\delta(\hat{t}), \\
0, & x \in \Omega \setminus U_\delta(\hat{t}), \end{cases}$$  

$$T^\lambda(x) = D\varepsilon(\nabla_x u^\lambda(x)), \quad x \in \Omega,$$  

where the function $\zeta \mapsto V(\lambda, \eta, \zeta) : [-\frac{\delta}{\lambda^{1/2}}, \frac{\delta}{\lambda^{1/2}}] \to \mathbb{R}^3$ is constructed as follows: We use the notations $V' = \partial_1 V$, $V'' = \partial_2^2 V$. In the interval $[-A(\mu), A(\mu)]$ the function $V$ is the solution of the boundary value problem

$$\left(D\varepsilon(V''(\lambda, \eta, \zeta) \otimes n)\right)n = \lambda^{1/2}f_{11}(\mu\lambda)(\eta, \lambda^{1/2}\zeta), \quad -A(\mu) \leq \zeta \leq A(\mu),$$  

$$V(\lambda, \eta, \pm A(\mu)) = 0,$$  

where $n = n(\eta)$ is the unit normal vector to $\Gamma(\hat{t})$ at $\eta \in \Gamma(\hat{t})$. The equation (6.24) is a second order linear system of ordinary differential equations for the three components of $V$, which can be written in the form

$$BV'' = \lambda^{1/2}f_{11}(\mu\lambda),$$  

(6.26)
with a $3 \times 3$-matrix $B = B(\eta)$ defined by the equation
\[ B\omega = (D\varepsilon(\omega \otimes n))n, \] (6.27)
which must hold for all $\omega \in \mathbb{R}^3$. The matrix $B$ is symmetric and positive definite uniformly with respect to $\eta$. To see this, note that since the elasticity tensor $D : S^3 \rightarrow S^3$ is a linear, symmetric, positive definite mapping, we compute for $\omega_1, \omega_2 \in \mathbb{R}^3$
\[ (B\omega_1) \cdot \omega_2 = \left( (D\varepsilon(\omega_1 \otimes n))n \right) \cdot \omega_2 
= (\omega_2 \otimes n) : D\varepsilon(\omega_1 \otimes n) = \varepsilon(\omega_2 \otimes n) : D\varepsilon(\omega_1 \otimes n) 
= (D\varepsilon(\omega_2 \otimes n)) : \varepsilon(\omega_1 \otimes n) = \left( (D\varepsilon(\omega_2 \otimes n))n \right) \cdot \omega_1 = (B\omega_2) \cdot \omega_1. \]

This shows that $B$ is symmetric. For $\omega \in \mathbb{R}^3$ we have with a suitable constant $C_0 > 0$, which only depends on $D$ but is independent of $\eta$, that
\[ (B\omega) \cdot \omega = \varepsilon(\omega \otimes n) : D\varepsilon(\omega \otimes n) \geq C_0|\varepsilon(\omega \otimes n)|^2 \geq \frac{C_0}{2}|\omega|^2, \]

hence $B$ is positive definite uniformly with respect to $\eta \in \Gamma(\hat{t})$.

Therefore the boundary value problem (6.24), (6.25) has a unique solution $V$ on $[-A(\mu), A(\mu)]$. To extend $\zeta \mapsto V(\lambda, \eta, \zeta)$ to all of $[-\frac{\delta}{\lambda^{1/2}}, \frac{\delta}{\lambda^{1/2}}]$, we continue $V$ to the intervals $(-\frac{\delta}{\lambda^{1/2}}, -A(\mu))$ and $(A(\mu), \frac{\delta}{\lambda^{1/2}})$ by affine functions:
\[ V(\lambda, \eta, \zeta) = \begin{cases} 
(\zeta + A(\mu))V'(\lambda, \eta, -A(\mu)), & -\frac{\delta}{\lambda^{1/2}} \leq \zeta \leq -A(\mu), \\
(\zeta - A(\mu))V'(\lambda, \eta, A(\mu)), & A(\mu) \leq \zeta \leq \frac{\delta}{\lambda^{1/2}}. 
\end{cases} \] (6.28)

By this extension, $\zeta \mapsto V(\lambda, \eta, \zeta)$ is continuously differentiable at $\zeta = \pm A(\mu)$. For $x = x(\hat{t}, \eta, \xi) \in U_{\delta}$ we use the notation
\[ V(\lambda, x) = V(\lambda, \eta, \frac{\xi}{\lambda^{1/2}}). \]

In the remaining part of the proof of Theorem 6.2 we need the following lemma, which we prove first.

**Lemma 6.3** There are constants $C_1, \ldots, C_4$ such that for all $\mu \in (0, \mu_0]$, $\lambda \in (0, \lambda_0]$, $(\eta, \zeta) \in \Gamma(\hat{t}) \times \left(-\frac{\delta}{\lambda^{1/2}}, \frac{\delta}{\lambda^{1/2}}\right)$ and $x \in U_{\delta}(\hat{t})$ the estimates
\[ |\nabla^j \xi V(\lambda, \eta, \zeta)| \leq C_1 |\ln \mu|^3 \mu, \quad j = 0, 1, 2, \] (6.29)
\[ |\lambda^{1/2} \nabla^j \xi V(\lambda, \eta, \zeta)| \leq C_2 |\ln \mu|^3 \mu, \quad j = 0, 1, 2, \] (6.30)
\[ |\lambda^{1/2} \nabla^j \xi V(\lambda, x)| \leq C_3 |\ln \mu|^3 \mu, \] (6.31)
\[ |\lambda^{1/2} \partial_{\xi_k} \nabla^j \xi V(\lambda, x)| \leq C_4 |\ln \mu|^3 \mu, \quad k = 1, \ldots, 3, \] (6.32)

hold. Moreover, there is a function $g^{(\mu)} : \Omega \rightarrow \mathbb{R}^3$ and a constant $C_5$ such that $T^{(\lambda)}$ defined in (6.29) satisfies
\[ \text{div}_x T^{(\lambda)} = f^{(\mu)} + g^{(\mu)}, \] (6.33)
with
\[ |g^{(\mu)}(x)| \leq C_5 |\ln \mu|^3 \mu, \] (6.34)

for all $x \in \Omega$ and all $\mu \in (0, \mu_0]$, $\lambda \in (0, \lambda_0]$.
\textbf{Proof:} In the following computations we drop the arguments \( \lambda \) and \( \eta \). Integration of (6.26) yields

\[
BV'(\zeta) = \int_{-A(\mu)}^{\zeta} \lambda^{1/2} f_{11}^{(\mu \lambda)}(\lambda^{1/2} \vartheta) d\vartheta + BV'(-A(\mu)),
\]

(6.35)

\[
BV(\zeta) = \int_{-A(\mu)}^{\zeta} \int_{-A(\mu)}^{\vartheta_1} \lambda^{1/2} f_{11}^{(\mu \lambda)}(\lambda^{1/2} \vartheta) d\vartheta d\vartheta_1 + (\zeta + A(\mu)) BV'(-A(\mu)),
\]

(6.36)

where we used the boundary condition (6.25) to get the second equation. Since \( V(A(\mu)) = 0 \), the relations (6.36) and (6.13) together yield

\[
2A(\mu) |BV'(A(\mu))| = \left| -\int_{-A(\mu)}^{A(\mu)} \int_{-A(\mu)}^{\vartheta_1} \lambda^{1/2} f_{11}^{(\mu \lambda)}(\lambda^{1/2} \vartheta) d\vartheta d\vartheta_1 \right|
\]

\[
\leq \int_{-A(\mu)}^{A(\mu)} \int_{-A(\mu)}^{\vartheta_1} |\ln \mu|^2 \mu^{1/2} K_1 d\vartheta d\vartheta_1 = 2A(\mu)^2 |\ln \mu|^2 \mu^{1/2} K_1,
\]

hence, by (6.10),

\[
|BV'(A(\mu))| \leq A(\mu) |\ln \mu|^2 \mu^{1/2} K_1 = \frac{2}{a} K_1 |\ln \mu|^3 \mu.
\]

Since \( B = B(\eta) \) is positive definite uniformly with respect to \( \eta \), this inequality implies the estimate (6.29) for \( j = 0 \) and \( -A(\mu) \leq \zeta \leq A(\mu) \). Since by definition (6.26) we have \( V'(\zeta) = V'(-A(\mu)) \) for \( \zeta \leq -A(\mu) \) and \( V'(\zeta) = V'(A(\mu)) \) for \( A(\mu) \leq \zeta \), the estimate (6.29) with \( j = 0 \) holds also for the values of \( \zeta \) outside of the interval \([-A(\mu), A(\mu)]\).

To prove (6.30) for \( j = 0 \) we use that \( V(-A(\mu)) = 0 \). By integration we thus obtain from (6.29) for \( \zeta \in \left[ \frac{-\delta}{V(\lambda)}, \frac{\delta}{V(\lambda)} \right] \) that

\[
|V(\zeta)| = \left| \int_{-A(\mu)}^{\zeta} V'(\vartheta) d\vartheta \right| \leq |\zeta + A(\mu)| C_1 |\ln \mu|^3 \mu \leq \left( \lambda^{-1/2} + A(\mu) \right) C_1 |\ln \mu|^3 \mu,
\]

which implies (6.30) for \( j = 0 \).

To verify (6.29) and (6.31) for \( j = 1, 2 \) we differentiate the differential equation (6.26) and the boundary condition (6.25) with respect to \( \eta \). For \( j = 1 \) we obtain the differential equation

\[
B(\eta)(\partial_{\eta} V)^n = \lambda^{1/2} (\partial_{\eta} f_{11}^{(\mu \lambda)} - \partial_{\eta} B(\eta) B(\eta)^{-1} f_{11}^{(\mu \lambda)}),
\]

and a similar equation for \( j = 2 \). We then use the estimate

\[
|\nabla_{\eta} j^{(\mu \lambda)}(\eta, \xi)| = |\nabla_{\eta} j^{(\mu \lambda)}(\eta, \xi)| \leq |\ln \mu|^2 \left( \frac{\mu}{\lambda} \right)^{1/2} K, \quad j = 1, 2.
\]

This estimate is obtained by differentiation with respect to \( \eta \) of the asymptotic expansions in Section 5.1 leading to Corollary 5.3. Under the regularity assumptions in Theorem 6.2 these derivatives exist. With this estimate we can employ the same arguments as above for the case \( j = 0 \) to derive (6.29) and (6.30) for \( j = 1, 2 \).

To prove (6.31) we use the decomposition (3.6) of the gradient and (3.13) to compute

\[
\nabla_{\lambda} \sum (\lambda, x) = \partial_{\xi} V(\lambda, \eta, \frac{\xi}{\lambda^{1/2}}) \otimes n + \nabla_{\Gamma} V(\lambda, \eta, \frac{\xi}{\lambda^{1/2}}) = \lambda^{-1/2} V' \otimes n + (\nabla_{\eta} V) A(i, \eta, \xi).
\]
The right hand side is estimated by (6.29) and (6.30) to obtain (6.31). The estimate (6.32) is obtained from (6.29) and (6.30) by similar decompositions.

To prove (6.33), note that by (6.22), (6.23) we have \( T^{(\lambda)} = D\varepsilon(\nabla_x u^{(\lambda)}) = D\varepsilon(\nabla_x (\lambda^{1/2} V \phi_*)) \). Using (3.6) and (3.7) we therefore obtain by a similar computation as in (5.73) that

\[
\text{div}_x T^{(\lambda)} = \text{div}_x D\varepsilon\left(\nabla_x (\lambda^{1/2} V (\lambda, \eta, \xi, \lambda^{1/2}) \phi_*)\right) \\
= \left(\lambda^{-1/2}(D\varepsilon(V'' \otimes n))n + \text{div}_{\Gamma}\varepsilon(D\varepsilon(V' \otimes n))\right)\phi_* \\
+ \lambda^{1/2}\left(D\varepsilon(\partial_t \nabla_{\Gamma} V))n + \text{div}_{\Gamma}\varepsilon(\nabla_{\Gamma} V)\right)\phi_* \\
+ \left((D\varepsilon(\lambda^{1/2} \nabla_x V))n + \text{div}_x D\varepsilon(\lambda^{1/2} V \otimes n))\right)\phi'_* \\
+ (D\varepsilon(\lambda^{1/2} V \otimes n))n \phi''_* \\
= f^{(\mu \lambda)}_1 + g^{(\mu \lambda)}.
\]

In the last step we used the differential equation (6.24) and noted that for \( \xi \in \left([-\delta, -A(\mu)\lambda^{1/2}] \cup [A(\mu)\lambda^{1/2}, \delta]\right) \) we have \( V''(\lambda, \eta, \xi, \lambda^{1/2}) = 0 \), by definition of \( V \) for such values of \( \xi \) in (6.28). We also used that \( \phi_* (\xi) = 1 \) for \( \xi \in [-A(\mu)\lambda^{1/2}, A(\mu)\lambda^{1/2}] \), which follows from (6.21) and (6.10), since we have chosen \( \mu_0 \) and \( \lambda_0 \) small enough such that \( A(\mu)\lambda^{1/2} < 2 \delta / 2 \) for all \( 0 < \mu \leq \mu_0 \) and \( 0 \leq \lambda \leq \lambda_0 \).

The function \( g^{(\mu \lambda)} \) is the sum of terms number 2 to 7 in the middle expression of equation (6.37). If we examine everyone of these six terms and apply (6.29) – (6.32) and also note that the functions \( \phi_* \), \( \phi'_* \) and \( \phi''_* \) are bounded independently of \( \mu \) and \( \lambda \) and vanish outside of \( \mathcal{U}_\delta(t) \), which follows from \( \phi_* \in C_0^\infty([-\delta, \delta]) \), we see that (6.34) holds for \( g^{(\mu \lambda)} \). This completes the proof of Lemma 6.3.

To conclude the proof of Theorem 6.2 let \((u^{(\lambda)}, T^{(\lambda)}_*)\) be the solution of the boundary value problem

\[
-\text{div}_x T^{(\lambda)}_* = g^{(\mu \lambda)} - f^{(\mu \lambda)}_1, \\
T^{(\lambda)}_* = D\varepsilon(\nabla_x u^{(\lambda)}), \\
u_1^{(\lambda)}(x) = 0, \quad x \in \partial\Omega.
\]

From these equations and from (6.22), (6.23), (6.33) we see that the function \((u^{(\lambda)} + u^{(\lambda)}_*, T^{(\lambda)} + T^{(\lambda)}_*)\) satisfies

\[
-\text{div}_x (T^{(\lambda)} + T^{(\lambda)}_*) = -f^{(\mu \lambda)}_1 - g^{(\mu \lambda)} + f^{(\mu \lambda)}_1 = -f^{(\mu \lambda)}_1(t), \\
(T^{(\lambda)} + T^{(\lambda)}_*) = D\varepsilon(\nabla_x (u^{(\lambda)} + u^{(\lambda)}_*)), \\
(u^{(\lambda)} + u^{(\lambda)}_*)(x) = 0, \quad x \in \partial\Omega,
\]

hence \((u^{(\lambda)} + u^{(\lambda)}_*, T^{(\lambda)} + T^{(\lambda)}_*)\) is equal to the unique solution of the boundary value problem (6.7) – (6.9), which means that \((u^{(\lambda)} + u^{(\lambda)}_*, T^{(\lambda)} + T^{(\lambda)}_*) = (u^{AC} - u, T^{AC} - u)(t)\).

Consequently, we have

\[
\|T^{AC} - T\|_{L^2(G(\hat{t}))} \leq \|T^{(\lambda)}\|_{L^2(G(\hat{t}))} + \|T^{(\lambda)}_*\|_{L^2(G(\hat{t}))}.
\]

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To estimate \( \| T^*(\lambda) \|_{L^2(\Gamma(\hat{t}))} \) we can use the theory of interior regularity for the elliptic boundary value problem (6.38) – (6.40). By this theory there is a constant \( C \) such that
\[
\| u^*(\lambda) \|_{W^{2,2}(\Omega')} \leq C \| g^{(\mu\lambda)} - f^{(\mu\lambda)}_{12} \|_{L^2(\Omega)},
\]
where \( \Omega' \) is the subdomain of \( \Omega \) introduced in Theorem 6.2, hence by the Sobolev embedding theorem and by (6.39),
\[
\| T^*(\lambda) \|_{L^2(\Gamma(\hat{t}))} \leq C_1 \| T^*(\lambda) \|_{W^{1,2}(\Omega')} \leq C_2 \| g^{(\mu\lambda)} - f^{(\mu\lambda)}_{12} \|_{L^2(\Omega')},
\]
where by our assumptions on \( \Gamma^{(\mu\lambda)}(\hat{t}) \) in Theorem 6.2 the constants \( C_1, C_2 \) can be chosen independently of \( \mu \) and \( \lambda \). By definition of \( f^{(\mu\lambda)}_{12} \) in (6.12) and by (2.52) we have
\[
| f^{(\mu\lambda)}_{12}(x) | \leq \mu \frac{3}{2} K^2 \text{ for all } x \in \Omega.
\]
From this inequality, from (6.34) and from (6.42) we conclude that
\[
\| T^*(\lambda) \|_{L^2(\Gamma(\hat{t}))} \leq C_2 \bigg( \int_{\Omega} (|g^{(\mu\lambda)}(x)| + |f^{(\mu\lambda)}_{12}(x)|)^2 \, dx \bigg)^{1/2}
\leq C_2 (C_5 | \ln \mu | \mu + \mu^{3/2} K_2) \bigg( \int_{\Omega} dx \bigg)^{1/2} \leq K | \ln \mu |^{3} \mu.
\]
From (6.23), (6.22) and from the inequalities (6.30), (6.31) we infer that
\[
| T^{(\lambda)}(x) | \leq C | \nabla_x u^{(\lambda)}(x) | = C | \nabla_x (\lambda^{1/2} V(\lambda, x) \phi^*_+(\xi)) | = C | (\lambda^{1/2} \nabla_x V(\lambda, x) \phi^*_+(\xi) + \lambda^{1/2} V(\lambda, x) \otimes (n \phi^*_+(\xi))) | \leq K' | \ln \mu |^3 \mu,
\]
whence
\[
\| T^{(\lambda)} \|_{L^2(\Gamma(\hat{t}))} \leq K'' | \ln \mu |^3 \mu.
\]
Combination of this inequality with (6.41) and (6.43) yields (6.4). The proof of Theorem 6.2 is complete.

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