# Optimal Control of time-discrete two-phase flow driven by a diffuse-interface model $\$$ 

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#### Abstract

We propose a general control framework for two-phase flows with variable densities in the diffuse interface formulation, where the distribution of the fluid components is described by a phase field. The flow is governed by the diffuse interface model proposed in [Abels, Garcke, Grün, M3AS $22(3): 1150013(40), 2012$ ]. On the basis of the stable time discretization proposed in [Garcke, Hinze, Kahle, APPL NUMER MATH, 99:151-171, 2016] we derive necessary optimality conditions for the time-discrete and the fully discrete optimal control problem. We present numerical examples with distributed and boundary controls, and also consider the case, where the initial value of the phase field serves as control variable. Keywords: Optimal control, Boundary control, Initial value control, Two-phase flow, Cahn-Hilliard, Navier-Stokes, Diffuse-interface models.


## 1 Introduction

In this paper we study a general discrete framework for control of two-phase fluids governed by the thermodynamically consistent diffuse interface model proposed in Abels et al., 2012. For the discretization we use the approach of Garcke et al., 2016|, where the authors propose a time discretization scheme, that preserves this important property in the time discrete setting and, using a post-processing step, also in the fully discrete setting including adaptive

[^0]mesh discretization. As control actions we consider distributed control, Dirichlet boundary control, and control with the initial condition of the phase field.

For the practical implementation we adapt the adaptive treatment developed in Garcke et al., 2016] to the optimal control setting. On the discrete level, special emphasis has to be be taken for the control with the intial value of the phase field, since the distribution of its phases is an outcome of the optimization procedure and thus a-priori unknown. In this case we combine the variational discretization from Hinze, 2005b with error estimation techniques to find a good mesh for the numerical representation of the a-priori unknown phase distribution.

Let us comment on related literature on time discretizations and control of (two-phase) fluids. For investigations of further time discretizations we refer to Aland, 2014, Hintermüller et al., 2015, Grün and Klingbeil, 2014, Garcke et al., 2016, Guillén-Gonzáles and Tierra, 2014, Guo et al., 2014 Grün et al., 2016. Concerning optimal control and feedback control of fluids there is a wide range of literature available. Here we only mention Gunzburger and Maservisi, 2000, Hinze and Kunisch, 2004, Fursikov et al., 1998, Berggren, 1998, Bewley et al., 2001. Hinze, 2005a].

Let us further comment on available literature for control of Cahn-Hilliard multiphase flow systems. In Hintermüller and Wegner, 2012 distributed optimal control of the Cahn-Hilliard system with a non smooth double obstacle potential is proposed, and in Hintermüller and Wegner, 2014] this work is extended to time-discrete two-phase flow given by a Cahn-Hilliard Navier-Stokes system with equal densities. Both works aim at existence of optimal controls and first order optimality conditions. In Hintermüller et al., 2015 the authors consider time discrete optimal control of multiphase flows based on the diffuse interface model of Abels et al., 2012. This work aims at establishing existence of solutions and stationarity conditions for control problems with free energies governed by the double obstacle potential, which is achieved through an appropriate limiting process of control problems with smooth relaxed free energies. The focus of the present work is different in that we consider numerical analysis of the fully discrete problem, propose a tailored numerical adaptive concept for the control problem, and present numerical examples which clearly show the potential of our approach.

We also mention the work of Baňas et al., 2014, where optimal control for a binary fluid, that is described by its density distribution, is proposed.

Let us finally comment on feedback control approaches for multiphase flows. Model predictive control is applied to the model from Abels et al., 2012 in Hinze and Kahle, 2013, Kahle, 2013. Kahle, 2014.

The paper is organized as follows. In Section 2 we state the model for the two-phase system and summarize assumptions that we require for the data. In Section 3 we state the time discretization scheme proposed in Garcke et al., 2016 and summarize properties of the scheme which we need in the present paper. We formulate the time discrete optimization problem in Section 3. In Section 4 we consider the optimal control problem in the fully discrete setting and present numerical examples in Section 5

## 2 The governing equations

The two-phase flow is modeled by the diffuse interface model proposed in Abels et al., 2012 .

$$
\begin{array}{rlrl}
\rho \partial_{t} v+((\rho v+J) \cdot \nabla) v-\operatorname{div}(2 \eta D v)+\nabla p & & \\
-\mu \nabla \varphi-\rho K-f & =0 & & \forall x \in \Omega, \forall t \in I, \\
-\operatorname{div}(v) & =0 & & \forall x \in \Omega, \forall t \in I, \\
\partial_{t} \varphi+v \cdot \nabla \varphi-\operatorname{div}(b \nabla \mu) & =0 & & \forall x \in \Omega, \forall t \in I, \\
-\sigma \epsilon \Delta \varphi+\frac{\sigma}{\epsilon} W^{\prime}(\varphi)-\mu & =0 & & \forall x \in \Omega, \forall t \in I, \\
v(0, x) & =v_{0}(x) & & \forall x \in \Omega, \\
\varphi(0, x) & =\varphi_{0}(x) & & \forall x \in \Omega, \\
v(t, x) & =g & & \forall x \in \partial \Omega, \forall t \in I, \\
\nabla \mu(t, x) \cdot \nu_{\Omega}=\nabla \varphi(t, x) \cdot \nu_{\Omega} & =0 & & \forall x \in \partial \Omega, \forall t \in I . \tag{8}
\end{array}
$$

Here $\varphi$ denotes the phase field, $\mu$ the chemical potential, $v$ the velocity field and $p$ the pressure. Furthermore $J=-\frac{\rho_{2}-\rho_{1}}{2} b \nabla \mu$ is a diffuse flux for $\varphi$.

In addition $\Omega \subset \mathbb{R}^{n}, n \in\{2,3\}$, denotes an open, convex and polygonal $(n=2)$ or polyhedral $(n=3)$ bounded domain. Its outer unit normal is denoted as $\nu_{\Omega}$, and $I=(0, T]$ with $0<T<\infty$ is a time interval.

The free energy density is denoted by $W$ and is assumed to be of double-well type with exactly two minima at $\pm 1$. For $W$ we use a splitting $W=W_{+}+W_{-}$, where $W_{+}$is convex and $W_{-}$is concave.

The density is denoted by $\rho=\rho(\varphi)$, fulfilling $\rho(-1)=\rho_{1}$ and $\rho(1)=\rho_{2}$, where $\rho_{1}, \rho_{2}$ denote the densities of the involved fluids. The viscosity is denoted by $\eta=\eta(\varphi)$, fulfilling $\eta(-1)=\eta_{1}$ and $\eta(1)=\eta_{2}$, with individual fluid viscosities $\eta_{1}, \eta_{2}$. The constant mobility is denoted by $b$. The gravitational force is denoted by $K$. By $D v=\frac{1}{2}\left(\nabla v+(\nabla v)^{t}\right)$ we denote the symmetrized gradient. The scaled surface tension is denoted by $\sigma$ and the interfacial width is proportional to $\epsilon$. We further have a volume force $f$ and boundary data $g$, as well as an initial phase field $\varphi_{0}$ and a solenoidal initial velocity field $v_{0}$.

Concerning results on existence of solutions for (1)-(8) under different assumptions on $W$ and $b$ we refer to Abels et al., 2013a, Abels et al., 2013b Grün, 2013.

## Assumptions

For the data of our problem we assume:
(A1) $W: \mathbb{R} \rightarrow \mathbb{R}$ is twice continuously differentiable and is of double-well type, i.e. it has exactly two minima at $\pm 1$ with values $W( \pm 1)=0$.
(A2) $W$ and its derivatives are polynomially bounded, i.e. there exists a $C>0$ such that $|W(x)| \leq C\left(1+|x|^{q}\right),\left|W_{+}^{\prime}(x)\right| \leq C\left(1+|x|^{q-1}\right),\left|W_{-}^{\prime}(x)\right| \leq C(1+$
$\left.|x|^{q-1}\right),\left|W_{+}^{\prime \prime}(x)\right| \leq C\left(1+|x|^{q-2}\right)$, and $\left|W_{-}^{\prime \prime}(x)\right| \leq C\left(1+|x|^{q-2}\right)$ holds for some $q \in[2,4]$ if $n=3$ and $q \in[2, \infty)$ if $n=2$.
(A3) There exists $\varphi_{a} \leq-1$ and $\varphi_{b} \geq 1$, such that $\rho(\varphi)=\rho\left(\varphi_{a}\right)$ for $\varphi \leq \varphi_{a}$, and $\rho(\varphi)=\rho\left(\varphi_{b}\right)$ for $\varphi \geq \varphi_{b}$. For $\varphi_{a}<\varphi<\varphi_{b}$ the function $\rho(\varphi)$ is affine linear, i.e. $\rho(\varphi)=\frac{1}{2}\left(\left(\rho_{2}-\rho_{1}\right) \varphi+\left(\rho_{1}+\rho_{2}\right)\right)$, and we define $\rho_{\delta}:=\frac{\left(\rho_{2}-\rho_{1}\right)}{2}$.

Further, $\eta(\varphi)=\eta\left(\varphi_{a}\right)$ for $\varphi \leq \varphi_{a}$, and $\eta(\varphi)=\eta\left(\varphi_{b}\right)$ for $\varphi \geq \varphi_{b}$. For $\varphi_{a}<$ $\varphi<\varphi_{b}$ the function $\eta(\varphi)$ is affine linear, i.e. $\eta(\varphi)=\frac{1}{2}\left(\left(\eta_{2}-\eta_{1}\right) \varphi+\left(\eta_{1}+\eta_{2}\right)\right)$. We define $\bar{\rho}>\underline{\rho}>0, \bar{\eta} \geq \underline{\eta}>0$ fulfilling

- $\bar{\rho} \geq \rho(\varphi) \geq \underline{\rho}>0$,
- $\bar{\eta} \geq \eta(\varphi) \geq \underline{\eta}>0$,
see Remark 2
(A4) The mean value of $\varphi$ is zero, i.e. there holds $\frac{1}{|\Omega|} \int_{\Omega} \varphi d x=0$. This can be achieved by choosing the values indicating the pure phases accordingly and considering a shifted system if required. In this case the values $\pm 1$ change to some other appropriate values.

Remark 1. The Assumptions (A1) (A2) are for example fulfilled by the polynomial free energy density

$$
W^{\text {poly }}(\varphi)=\frac{1}{4}\left(1-\varphi^{2}\right)^{2} .
$$

Another free energy density fulfilling these assumptions is the relaxed doubleobstacle free energy density given by

$$
\begin{align*}
\lambda(y) & :=\max (0, y-1)+\min (0, y+1), \\
\xi & :=\frac{1+2 s+\sqrt{4 s+1}}{2 s}, \\
\delta & :=\frac{1}{2}\left(1-\xi^{2}\right)+\frac{s}{3}|\lambda(\xi)|^{3}, \\
W^{s}(y) & =\frac{1}{2}\left(1-(\xi y)^{2}\right)+\frac{s}{3}|\lambda(\xi y)|^{3}+\delta \tag{9}
\end{align*}
$$

where $s \gg 0$ denotes a relaxation parameter. $W^{s}$ can be understood as a relaxation of the double-obstacle free energy density

$$
W^{\infty}(\varphi)= \begin{cases}\frac{1}{2}\left(1-\varphi^{2}\right) & \text { if }|\varphi| \leq 1 \\ 0 & \text { else }\end{cases}
$$

which is proposed in Oono and Puri, 1988, Blowey and Elliott, 1991 to model phase separation. We note that here we use a cubic penalisiation to obtain the required regularity from (A1) and that $\xi$ is chosen such that $W^{s}$ takes its minima at $\pm 1$ and $\delta$ is such that $W^{s}( \pm 1)=0$.

In the numerical examples of this work we use the free energy density $W \equiv$ $W^{s}$. For this choice the splitting into convex and concave part reads

$$
W_{+}(\varphi)=s \frac{1}{3}|\lambda(\xi \varphi)|^{3}, \quad W_{-}(\varphi)=\frac{1}{2}\left(1-(\xi \varphi)^{2}\right)+\delta
$$

Remark 2. For the weak formulation of (1)-(8) we later require affine linearity of $\rho$ on the image of $\varphi$. The affine linearity of $\eta$ is assumed for simplicity. Note that in view of Assumption (A3), this essentially implies a bound on $\varphi$, namely $\varphi \in\left(\varphi_{a}, \varphi_{b}\right)$ as stated in Assumption (A3).

Using $W^{s}$ as free energy density we argue, that for sufficiently large (see |Garcke et al., 2016, Rem. 6]) $|\varphi| \leq 1+\theta$ holds, with $\theta$ sufficiently small, and in [Kahle, 2015] it is shown for the Cahn-Hilliard equation without transport, that for the energy (9) in fact $\|\varphi\|_{L^{\infty}(\Omega)} \leq 1+C s^{-1 / 2}$ holds.

In a general setting one might use a nonlinear dependence between $\varphi$ and $\rho$, see e.g. [Abels and Breit, 2016], or choose a cut-off procedure as proposed in Grün, 2013, Guillén-Gonzáles and Tierra, 2014].

Anyway, since we later require linearity of $\rho$ on the image of $\varphi$ we state Assumption (A3) and note that this assumption is fulfilled in our numerical examples in Section 5.

## Notation

We use the conventional notation for Sobolev and Hilbert Spaces, see e.g. Adams and Fournier, 2003. With $L^{p}(\Omega), 1 \leq p \leq \infty$, we denote the space of measurable functions on $\Omega$, whose modulus to the power $p$ is Lebesgueintegrable. $L^{\infty}(\Omega)$ denotes the space of measurable functions on $\Omega$, which are essentially bounded. For $p=2$ we denote by $L^{2}(\Omega)$ the space of square integrable functions on $\Omega$ with inner product $(\cdot, \cdot)$ and norm $\|\cdot\|$. By $W^{k, p}(\Omega)$, $k \geq 1,1 \leq p \leq \infty$, we denote the Sobolev space of functions admitting weak derivatives up to order $k$ in $L^{p}(\Omega)$. If $p=2$ we write $H^{k}(\Omega)$.

For $f \in H^{1}(\Omega)^{n}$ we introduce the continuous trace operator $\gamma: H^{1}(\Omega)^{n} \rightarrow$ $H^{\frac{1}{2}}(\partial \Omega)^{n}$ as $\gamma f:=\left.f\right|_{\partial \Omega}$. We further note that for $g \in H^{\frac{1}{2}}(\partial \Omega)^{n}$ with $g \cdot \nu_{\Omega}=0$ there exists $\widetilde{g} \in H^{1}(\Omega)^{n},(\operatorname{div} \widetilde{g}, q)=0 \forall q \in L^{2}(\Omega)$ with $\gamma \widetilde{g}=g$ and $\|\widetilde{g}\|_{H^{1}(\Omega)^{n}} \leq$ $C\|g\|_{H^{\frac{1}{2}}(\partial \Omega)^{n}}$, where $C$ is independent of $g$.

The subspace $H_{0}^{1}(\Omega)^{n} \subset H^{1}(\Omega)^{n}$ denotes the set of functions with vanishing boundary trace. We further set

$$
L_{(0)}^{2}(\Omega)=\left\{v \in L^{2}(\Omega) \mid(v, 1)=0\right\}
$$

and with

$$
H_{\sigma}(\Omega)=\left\{v \in H^{1}(\Omega)^{n} \mid(\operatorname{div}(v), q)=0 \forall q \in L^{2}(\Omega)\right\}
$$

we denote the space of all weakly solenoidal $H^{1}(\Omega)$ vector fields. We stress that there is no correspondence between the subscript $\sigma$ and the scaled surface
tension. We denote both terms using $\sigma$ since these are standard notations. We further introduce

$$
H_{0, \sigma}(\Omega)=H_{0}^{1}(\Omega)^{n} \cap H_{\sigma}(\Omega)
$$

For $u \in L^{q}(\Omega)^{n}, q>2$ if $n=2, q \geq 3$ if $n=3$, and $v, w \in H^{1}(\Omega)^{n}$ we introduce the trilinear form

$$
\begin{equation*}
a(u, v, w)=\frac{1}{2} \int_{\Omega}((u \cdot \nabla) v) w d x-\frac{1}{2} \int_{\Omega}((u \cdot \nabla) w) v d x \tag{10}
\end{equation*}
$$

Note that there holds $a(u, v, w)=-a(u, w, v)$, and especially $a(u, v, v)=0$. We have the following stability estimate by Hölder inequalities and Sobolev embedding

$$
|a(u, v, w)| \leq C\|u\|_{L^{q}(\Omega)}\|v\|_{H^{1}(\Omega)}\|w\|_{H^{1}(\Omega)}
$$

For a square summable series of functions $\left(f_{m}\right)_{m=1}^{M} \in V^{M}$, where $\left(V,\|\cdot\|_{V}\right)$ is a normed vector space, we introduce the notation $\left\|\left(f^{m}\right)_{m=1}^{M}\right\|_{V}^{2}=\sum_{m=1}^{M}\left\|f_{m}\right\|_{V}^{2}$.

## 3 The time-discrete setting

In Garcke et al., 2016 existence of time discrete weak solutions for (1)-(4) is shown for the case of $g=0$ and $f=0$. In this section we formulate a time discrete optimization problem for (11)- (4), where we use $g, f$, and $\varphi^{0}$ as controls, and show existence of solutions together with first order optimality conditions.

Let $0=t_{0}<t_{1}<\ldots<t_{m-1}<t_{m}<t_{m+1}<\ldots<t_{M}=T$ denote an equidistant subdivision of the interval $\bar{I}=[0, T]$ with $\tau_{m+1}-\tau_{m} \equiv \tau$ and sub intervals $I_{0}=\{0\}, I_{m}=\left(t_{m-1}, t_{m}\right], m=1, \ldots, M$. From here onwards the superscript $m$ denotes the corresponding variables at time instance $t_{m}$, e.g. $\varphi^{m}:=\varphi\left(t_{m}\right)$. For functions $f \in L^{2}(0, T, V)$ we introduce $f^{m}:=f_{I_{m}} f(t) d t \in V$. Note that this can be seen as a discontinuous Galerkin approximation using piecewise constant values.

We now introduce the optimal control problem under consideration. For this purpose we interpret $\varphi_{0}, f$, and $g$ as sought control that we intend to choose, such that the corresponding phase field $\varphi^{M}$ is close to a desired phase field $\varphi_{d}$ in the mean square sense. If $\varphi_{d}$ is the measurement of a real world system, then finding $\varphi^{0}$ such that the corresponding phase field $\varphi^{M}$ is close to $\varphi_{d}$ resembles an inverse problem.

We denote by $u \in U$ the control, where

$$
U=U_{I} \times U_{V} \times U_{B}=\mathcal{K} \times L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right) \times L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)
$$

is the space of controls, where

$$
\mathcal{K}:=\left\{v \in H^{1}(\Omega)\left|\int_{\Omega} v d x=0,|v| \leq 1\right\} \subset H^{1}(\Omega) \cap L^{\infty}(\Omega)\right.
$$

denotes the space of admissible initial phase fields.
By

$$
\mathcal{B}: U \rightarrow H^{1}(\Omega) \cap L^{\infty}(\Omega) \times L^{2}\left(0, T ; L^{2}(\Omega)^{n}\right) \times L^{2}\left(0, T ;\left(H^{1 / 2}(\partial \Omega)\right)^{n}\right)
$$

we denote the linear and bounded control operator, which consists of three components, i.e. $\mathcal{B}=\left[B_{I}, B_{V}, B_{B}\right]$, where $B_{I}\left(u_{I}, u_{V}, u_{B}\right) \equiv B_{I} u_{I}:=u_{I}$, which is the initial phase field for the system, $B_{V}\left(u_{I}, u_{V}, u_{B}\right) \equiv B_{V} u_{V}$ with $B_{V} u_{V}(t, x)=$ $\sum_{l=1}^{u_{v}} f_{l}(x) u_{V}^{l}(t)$ where $f_{l} \in L^{2}(\Omega)^{n}$ are given functions, which is a volume force acting on the fluid inside $\Omega$, and $B_{B}\left(u_{I}, u_{V}, u_{B}\right) \equiv B_{B} u_{B}$, with $B_{B} u_{B}(t, x)=$ $\sum_{l=1}^{u_{b}} g_{l}(x) u_{B}^{l}(t)$ where $g_{l} \in H^{1 / 2}(\partial \Omega)^{n}$ denote given functions, and this is a boundary force acting on the fluid as Dirichlet boundary data. To obtain a solenoidal velocity field, $B_{B} u_{B}$ has to fulfill the compatibility condition $\int_{\partial \Omega} B_{B} u_{B}$. $\nu_{\Omega} d s=0$, and in the following for simplicity we assume $g_{l} \cdot \nu_{\Omega}=0, l=1, \ldots, u_{b}$, point wise.

Given a triple $\left(\alpha_{I}, \alpha_{V}, \alpha_{B}\right)$ of non negative values with $\alpha_{I}+\alpha_{V}+\alpha_{B}=1$ we introduce an inner product for $u=\left(u_{I}, u_{V}, u_{B}\right) \in U$ and $v=\left(v_{I}, v_{V}, v_{B}\right) \in U$ by

$$
\begin{equation*}
(u, v)_{U}=\alpha_{I}\left(\nabla u_{I}, \nabla v_{I}\right)_{L^{2}(\Omega)}+\alpha_{V}\left(u_{V}, v_{V}\right)_{L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right)}+\alpha_{B}\left(u_{B}, v_{B}\right)_{L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)} \tag{11}
\end{equation*}
$$

and the norm $\|u\|_{U}^{2}=(u, u)_{U}$.
We use the convention, that $\alpha_{\star}=0, \star \in\{I, V, B\}$, means, that we do not apply this kind of control. If $\alpha_{I}=0$ we use $\varphi^{0}$ as given data, if $\alpha_{B}=0$, we assume no-slip boundary data for $v$. For notational convenience, in the following we assume $\alpha_{\star} \neq 0$ for all $\star \in\{I, V, B\}$.

We stress, that we do not discretize the control in time, although the state equation is time discrete. Thus we follow the concept of variational discretization Hinze, 2005b. Anyway, the control is discretized implicitly in time by the adjoint equation that we will derive later. We also note, that in view of the state equation, this allows us to dynamically adapt the time step size $\tau$ to the flow condition without changing the control space.

Following Garcke et al., 2016 we propose the following time discrete counterpart of (1)-(8):
Let $u \in U$ and $v_{0} \in H_{\sigma}(\Omega) \cap L^{\infty}(\Omega)$ be given.
Initialization for $m=1$ :
Set $\varphi^{0}=u_{I}$ and $v^{0}=v_{0}$.
Find $\varphi^{1} \in H^{1}(\Omega) \cap L^{\infty}(\Omega), \mu^{1} \in W^{1,3}(\Omega), v^{1} \in H_{\sigma}(\Omega)$, with $\gamma\left(v^{1}\right)=B_{B} u_{B}^{1}$, such
that for all $w \in H_{0, \sigma}(\Omega), \Phi \in H^{1}(\Omega)$, and $\Psi \in H^{1}(\Omega)$ it holds

$$
\begin{array}{r}
\frac{1}{\tau} \int_{\Omega}\left(\frac{1}{2}\left(\rho^{1}+\rho^{0}\right) v^{1}-\rho^{0} v^{0}\right) w d x+a\left(\rho^{1} v^{0}+J^{1}, v^{1}, w\right) \\
+\int_{\Omega} 2 \eta^{1} D v^{1}: D w d x-\int_{\Omega} \mu^{1} \nabla \varphi^{0} w+\rho^{0} K w d x-\left\langle B_{V} u_{V}^{1}, w\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}=0 \\
\frac{1}{\tau} \int_{\Omega}\left(\varphi^{1}-\varphi^{0}\right) \Psi d x+\int_{\Omega}\left(v^{0} \cdot \nabla \varphi^{0}\right) \Psi d x+\int_{\Omega} b \nabla \mu^{1} \cdot \nabla \Psi d x=0 \\
\sigma \epsilon \int_{\Omega} \nabla \varphi^{1} \cdot \nabla \Phi d x-\int_{\Omega} \mu^{1} \Phi d x+\frac{\sigma}{\epsilon} \int_{\Omega}\left(W_{+}^{\prime}\left(\varphi^{1}\right)+W_{-}^{\prime}\left(\varphi^{0}\right)\right) \Phi d x=0 \tag{14}
\end{array}
$$

where $J^{1}:=-\rho_{\delta} b \nabla \mu^{1}$.
Two-step scheme for $m>1$ :
Given $\varphi^{m-2} \in H^{1}(\Omega) \cap L^{\infty}(\Omega), \varphi^{m-1} \in H^{1}(\Omega) \cap L^{\infty}(\Omega), \mu^{m-1} \in W^{1,3}(\Omega), v^{m-1} \in$ $H_{\sigma}(\Omega)$,
find $v^{m} \in H_{\sigma}(\Omega), \gamma\left(v^{m}\right)=B_{B} u_{B}^{m}, \varphi^{m} \in H^{1}(\Omega) \cap L^{\infty}(\Omega), \mu^{m} \in W^{1,3}(\Omega)$ such that for all $w \in H_{0, \sigma}(\Omega), \Psi \in H^{1}(\Omega)$, and $\Phi \in H^{1}(\Omega)$ it holds

$$
\begin{array}{r}
\frac{1}{\tau} \int_{\Omega}\left(\frac{\rho^{m-1}+\rho^{m-2}}{2} v^{m}-\rho^{m-2} v^{m-1}\right) w d x+\int_{\Omega} 2 \eta^{m-1} D v^{m}: D w d x \\
+a\left(\rho^{m-1} v^{m-1}+J^{m-1}, v^{m}, w\right) \\
-\int_{\Omega} \mu^{m} \nabla \varphi^{m-1} w+\rho^{m-1} K w d x-\left\langle B_{V} u_{V}^{m}, w\right\rangle_{H^{-1}(\Omega), H_{0}^{1}(\Omega)}=0 \\
\int_{\Omega} \frac{\varphi^{m}-\varphi^{m-1}}{\tau} \Psi d x+\int_{\Omega}\left(v^{m} \cdot \nabla \varphi^{m-1}\right) \Psi d x+\int_{\Omega} b \nabla \mu^{m} \cdot \nabla \Psi d x=0 \\
\sigma \epsilon \int_{\Omega} \nabla \varphi^{m} \cdot \nabla \Phi d x-\int_{\Omega} \mu^{m} \Phi d x+\frac{\sigma}{\epsilon} \int_{\Omega}\left(W_{+}^{\prime}\left(\varphi^{m}\right)+W_{-}^{\prime}\left(\varphi^{m-1}\right)\right) \Phi d x=0 \tag{17}
\end{array}
$$

where $J^{m-1}:=-\rho_{\delta} b \nabla \mu^{m-1}$. We further use the abbreviations $\rho^{m}:=\rho\left(\varphi^{m}\right)$ and $\eta^{m}:=\eta\left(\varphi^{m}\right)$.

We note that in 15 - 17 ) the only nonlinearity arises from $W_{+}^{\prime}$ and thus only the equation $\sqrt{17 /}$ is nonlinear. A similar argumentation holds for $\sqrt{12}$ (14). The regularity $\nabla \mu^{m-1} \in L^{3}(\Omega)$ is required for the trilinear form $a(\cdot, \cdot, \cdot)$, see (10).

Remark 3. We note that (15)-(17) is a two-step scheme for the phase field variable $\varphi$, and thus we need an initialization as proposed in 12- 14). Here, as in [Garcke et al., 2016] the sequential coupling of (13)-14] and 12 is used as proposed in [Kay et al., 2008].

Another variant might be to require initial data on time instance $t_{-1}$ for the phase field and at $t_{0}$ for the velocity field. Equations (16) -(17) can than be solved for $\varphi^{0}$ and $\mu^{0}$ to obtain initial values, see Hintermüller et al., 2015.

Since we are later also interested in control of the initial value $\varphi_{0}$ we propose the initialization scheme $\sqrt{12}-(\sqrt{14})$ here.

Theorem 4. Let $v^{0} \in H_{\sigma}(\Omega) \cap L^{\infty}(\Omega)^{n}$ and $u \in U$ be given data.
Then there exists a unique solution $\left(v^{1}, \varphi^{1}, \mu^{1}\right)$ to (12)-14), and it holds

$$
\begin{align*}
&\left\|v^{1}\right\|_{H^{1}(\Omega)^{n}}+\left\|\varphi^{1}\right\|_{H^{2}(\Omega)}+\left\|\mu^{1}\right\|_{H^{2}(\Omega)} \\
& \leq C_{1}\left(v_{0}\right) C_{2}\left(\left\|u_{I}\right\|_{H^{1}(\Omega)},\left\|B_{V} u_{V}^{1}\right\|_{L^{2}(\Omega)^{n}},\left\|B_{B} u_{B}^{1}\right\|_{H^{\frac{1}{2}}(\partial \Omega)^{n}}\right) \tag{18}
\end{align*}
$$

and $\varphi^{1}, \mu^{1}$ can be found be Newton's method. The constant $C_{2}$ depends polynomially on its arguments.

Proof. The existence of $\left(\varphi^{1}, \mu^{1}\right) \in H^{1}(\Omega) \times H^{1}(\Omega)$ follows from (Hintermüller et al., 2011) $)$. There the corresponding system without the transport term $v^{0} \nabla u_{I}$ is analyzed. This term is a given volume force, that can be incorporated in a straightforward manner. From this we directly obtain the stability inequality

$$
\left\|\varphi^{1}\right\|_{H^{1}(\Omega)}+\left\|\mu^{1}\right\|_{H^{1}(\Omega)} \leq C_{1}\left(v^{0}\right) C_{2}\left(\left\|u_{I}\right\|_{H^{1}(\Omega)}\right)
$$

Since $\left|W_{+}^{\prime}(\varphi)\right| \leq C\left(1+|\varphi|^{q} \mid\right), q \leq 3$ we have $W_{+}^{\prime}(\varphi) \in L^{2}(\Omega)$ and by $L^{2}$ regularity theory we have $\varphi^{1} \in H^{2}(\Omega)$ and

$$
\left\|\varphi^{1}\right\|_{H^{2}(\Omega)} \leq C\left(\left\|\mu^{1}\right\|_{H^{1}(\Omega)},\left\|\varphi^{1}\right\|_{H^{1}(\Omega)},\left\|u_{I}\right\|_{H^{1}(\Omega)}\right)
$$

We further have $v^{0} \nabla u_{I} \in L^{2}(\Omega)$ and thus we have $\mu^{1} \in H^{2}(\Omega)$ and the stability inequality

$$
\left\|\mu^{1}\right\|_{H^{2}(\Omega)} \leq C_{1}\left(v_{0}\right) C_{2}\left(\left\|u_{I}\right\|_{H^{1}(\Omega)},\left\|\varphi^{1}\right\|_{H^{1}(\Omega)}\right)
$$

Convergence of Newton's method directly follows from Hintermüller et al., 2011. Note that the only nonlinearity $W_{+}^{\prime}$ is monotone.

With $v^{0}, \varphi^{1}, u_{I}$, and $\mu^{1}$ given data, (12) defines a coercive and continuous bilinear form on $H_{\sigma}$ and thus existence and stability of a solution follows from Lax-Milgram's theorem. This uses the antisymmetry of the trilinear form $a$ and Korn's inequality.

Theorem 5. Let $v^{m-1} \in H_{\sigma}(\Omega), \varphi^{m-2} \in H^{1}(\Omega) \cap L^{\infty}(\Omega), \varphi^{m-1} \in H^{1}(\Omega) \cap$ $L^{\infty}(\Omega)$, and $\mu^{m-1} \in W^{1,3}(\Omega)$, be given data. Then there exists a unique solution $\left(v^{m}, \varphi^{m}, \mu^{m}\right)$ to (15)-17).

It further holds $\varphi^{m} \in H^{2}(\Omega)$ and if additionally $\varphi^{m-1} \in W^{1,3}(\Omega)$ we have $\mu^{m} \in H^{2}(\Omega)$ and the stability inequality

$$
\begin{aligned}
& \left\|v^{m}\right\|_{H^{1}(\Omega)^{n}}+\left\|\mu^{m}\right\|_{H^{2}(\Omega)}+\left\|\varphi^{m}\right\|_{H^{2}(\Omega)} \\
& \leq C\left(\left\|v^{m-1}\right\|_{H^{1}(\Omega)^{n}},\left\|\varphi^{m-1}\right\|_{W^{1,3}(\Omega)},\left\|\varphi^{m-2}\right\|_{H^{1}(\Omega)},\right. \\
& \left.\quad\left\|B_{V} u_{V}^{m}\right\|_{L^{2}(\Omega)^{n}},\left\|B_{B} u_{B}^{m}\right\|_{H^{\frac{1}{2}}(\partial \Omega)^{n}}\right)
\end{aligned}
$$

holds. The constant $C$ depends polynomially on its arguments.

Proof. In Garcke et al., 2016 the existence for $\gamma\left(B_{B} u_{B}^{m}\right)=0$ and $B_{V} u_{V}^{m}=0$ is shown using a Galerkin approach. The additional volume force is incorporated in a straight forward manner, and the boundary data $B_{B} u_{B}^{m}$ can be introduced by investigating a shifted system, see Theorem 4

We define $e:=\widehat{B_{B} u_{B}^{m}}$ and use $w=v^{m}-e$ as test function in 15, $\Psi=\mu^{m}$ as test function in 16), and $\Phi=\tau^{-1}\left(\varphi^{m}-\varphi^{m-1}\right)$ as test function in 17), and add the resulting equations. Using the properties of $W_{+}^{\prime}$ and $W_{-}^{\prime}$ we obtain (compare Garcke et al., 2016, Thm. 3])

$$
\begin{aligned}
& E\left(v^{m}, \varphi^{m}, \varphi^{m-1}\right)+\frac{1}{2} \int_{\Omega} \rho^{m-2}\left|v^{m}-v^{m-1}\right|^{2} d x+2 \tau \int_{\Omega} \eta^{m-1}\left|D v^{m}\right|^{2} d x \\
& +\tau \int_{\Omega} b\left|\nabla \mu^{m}\right|^{2} d x+\frac{\sigma \epsilon}{2}\left\|\nabla \varphi^{m}-\nabla \varphi^{m-1}\right\|^{2} \\
& \leq E\left(v^{m-1}, \varphi^{m-1}, \varphi^{m-2}\right)+\int_{\Omega}\left(\frac{\rho^{m-1}+\rho^{m-2}}{2} v^{m}-\rho^{m-2} v^{m-1}\right) e d x \\
& +\tau a\left(\rho^{m-1} v^{m-1}+J^{m-1}, v^{m}, e\right)+2 \tau \int_{\Omega} \eta^{m-1} D v^{m}: D e d x \\
& -\tau \int_{\Omega} \mu^{m} \nabla \varphi^{m-1} e d x+\left(\rho^{m-1} K, v^{m}-e\right)+\tau\left(B_{V} u_{V}^{m}, v^{m}-e\right)_{L^{2}(\Omega)^{n}}
\end{aligned}
$$

By using the inequalities of Hölder, Korn and Young, together with Assumption (A3) and the stability of the extension operator $\sim$ the claim follows. The regularity $\varphi^{m}, \mu^{m} \in H^{2}(\Omega)$ follow as in the proof of Theorem 4, but now using $\nabla \varphi^{m-1} \in L^{3}(\Omega)$ and $v^{m} \in H_{\sigma} \hookrightarrow L^{6}(\Omega)$.

Let us next introduce the optimization problem under investigation. For this we first rewrite $12-17$ in a compact and abstract form and introduce

$$
\begin{gather*}
Y:=H_{\sigma}(\Omega)^{M} \times\left(H^{1}(\Omega) \cap L^{\infty}(\Omega)\right)^{M} \times W^{1,3}(\Omega)^{M}, \\
Y_{0}:=H_{0, \sigma}(\Omega)^{M} \times\left(H^{1}(\Omega) \cap L^{\infty}(\Omega)\right)^{M} \times W^{1,3}(\Omega)^{M}, \\
y:=\left(v^{m}, \varphi^{m}, \mu^{m}\right)_{m=1}^{M} \in Y, \\
Z:=\left(H_{0, \sigma}(\Omega)^{M} \times H^{1}(\Omega)^{M} \times H^{1}(\Omega)^{M}\right)^{\star} \\
e: Y_{0} \times U \rightarrow Z, \\
e\left(y_{0}, u\right)=0 \tag{19}
\end{gather*}
$$

The operator $e$ is defined as follows

$$
\begin{aligned}
& \left\langle\tilde{y}, e\left(y_{0}, u\right)\right\rangle_{Z^{\star}, Z}:= \\
& \tau^{-1}\left(\frac{1}{2}\left(\rho^{1}+\rho^{0}\right)\left(v_{0}^{1}+\widetilde{B_{B} u_{B}^{1}}\right)-\rho^{0} v^{0}, \tilde{v}^{1}\right)+a\left(\rho^{1} v^{0}+J^{1}, v_{0}^{1}+\widetilde{B_{B} u_{B}^{1}}, \tilde{v}^{1}\right) \\
& +\left(2 \eta^{1} D\left(v_{0}^{1}+\widetilde{B_{B} u_{B}^{1}}\right), D \tilde{v}^{1}\right)-\left(\mu^{1} \nabla u_{1}+\rho^{0} K, \tilde{v}^{1}\right) \\
& -\left(B_{V} u_{V}^{1}, \tilde{v}^{1}\right) \\
& +\tau^{-1}\left(\varphi^{1}-u_{1}, \tilde{\varphi}^{1}\right)+\left(v^{0} \nabla u_{1}, \tilde{\varphi}^{1}\right)+\left(b \nabla \mu^{1}, \nabla \tilde{\varphi}^{1}\right) \\
& +\sigma \epsilon\left(\nabla \varphi^{1}, \nabla \tilde{\mu}^{1}\right)-\left(\mu^{1}, \tilde{\mu}^{1}\right) \\
& +\sigma \epsilon^{-1}\left(W_{+}^{\prime}\left(\varphi^{1}\right)+W_{-}^{\prime}\left(u_{1}\right), \tilde{\mu}^{1}\right) \\
+\sum_{m=2}^{M} & {\left[\tau^{-1}\left(\frac{1}{2}\left(\rho^{m-1}+\rho^{m-2}\right)\left(v_{0}^{m}+\widetilde{B_{B} u_{B}^{m}}\right)-\rho^{m-2}\left(v_{0}^{m-1}+\widetilde{B_{B} u_{B}^{m-1}}\right), \tilde{v}^{m}\right)\right.} \\
& +a\left(\rho^{m-1}\left(v_{0}^{m-1}+\widehat{B_{B} u_{B}^{m-1}}\right)+J^{m-1}, v_{0}^{m}+\widetilde{B_{B} u_{B}^{m}}, \tilde{v}^{m}\right) \\
& +\left(2 \eta^{m-1} D\left(v_{0}^{m}+\widehat{B_{B} u_{B}^{m}}\right), D \tilde{v}^{m}\right)-\left(\mu^{m} \nabla \varphi^{m-1}+\rho^{m-1} K, \tilde{v}^{m}\right) \\
& -\left(B_{V} u_{V}^{m}, \tilde{v}^{m}\right) \\
& +\tau^{-1}\left(\varphi^{m}-\varphi^{m-1}, \tilde{\varphi}^{m}\right)+\left(\left(v_{0}^{m}+\widetilde{B_{B} u_{B}^{m}}\right) \nabla \varphi^{m-1}, \tilde{\varphi}^{m}\right)+\left(b \nabla \mu^{m}, \nabla \tilde{\varphi}^{m}\right) \\
& +\sigma \epsilon\left(\nabla \varphi^{m}, \nabla \tilde{\mu}^{m}\right)-\left(\mu^{m}, \tilde{\mu}^{m}\right) \\
& \left.+\sigma \epsilon^{-1}\left(W_{+}^{\prime}\left(\varphi^{m}\right)+W_{-}^{\prime}\left(\varphi^{m-1}\right), \tilde{\mu}^{m}\right)\right]
\end{aligned}
$$

with $y_{0}:=\left(v_{0}^{m}, \varphi^{m}, \mu^{m}\right)_{m=1}^{M} \in Y_{0}$, and $\left.\left.\left.\tilde{y}=\left(\left(\tilde{v}^{m}\right)_{m=1}^{M},\left(\tilde{\varphi}^{m}\right)_{m=1}^{M},\left(\tilde{\mu}^{m}\right)_{m=1}^{M}\right)\right)\right)\right) \epsilon$ $Z^{\star}$. Here again $\rho^{m}:=\rho\left(\varphi^{m}\right), \eta^{m}:=\eta\left(\varphi^{m}\right)$ and especially $\rho^{0}:=\rho\left(u_{I}\right), \eta^{0}:=$ $\eta\left(u_{I}\right)$.

Now the time-discrete optimization problem under investigation is given as

$$
\begin{align*}
\min _{u \in U} J\left(\left(\varphi^{m}\right)_{m=1}^{M}, u\right)=\frac{1}{2} \| & \varphi^{M}-\varphi_{d} \|_{L^{2}(\Omega)}^{2} \\
& +\frac{\alpha}{2}\left(\alpha_{I} \int_{\Omega} \frac{\epsilon}{2}\left|\nabla u_{I}\right|^{2}+\epsilon^{-1} W_{u}\left(u_{I}\right) d x\right.  \tag{P}\\
& \left.+\alpha_{V}\left\|u_{V}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right)}^{2}+\alpha_{B}\left\|u_{B}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)}^{2}\right)
\end{align*}
$$

s.t. $e(y, u)=0$.

Here $\varphi_{d} \in L^{2}(\Omega)$ is a given desired phase field, and $\alpha>0$ is a weight for the control cost. For the control cost of the initial value we use the well-known Ginzburg-Landau energy of the phase field $u_{I}$ with interfacial thickness $\epsilon$. Here we use the double obstacle free energy density $W_{u} \equiv W^{\infty}$ given in Remark 1. In our numerical examples it is advantageous to use this non-smooth free energy density instead of the smoother one used for the simulation.

Theorem 6. Let $v^{0} \in H_{\sigma}(\Omega) \cap L^{\infty}(\Omega), u \in U$ be given.

Then there exists a unique solution to the equation $e(y, u)=0$, i.e. there exist $\left(v^{m}, \varphi^{m}, \mu^{m}\right)_{m=1}^{M} \in Y$ such that $\left(v^{m}, \varphi^{m}, \mu^{m}\right)$ is the unique solution to $12-17$ ) for $m=1, \ldots, M$. Moreover there holds

$$
\begin{aligned}
& \left\|\left(v^{m}\right)_{m=1}^{M}\right\|_{H^{1}(\Omega)^{n}}+\left\|\left(\varphi^{m}\right)_{m=1}^{M}\right\|_{H^{2}(\Omega)}+\left\|\left(\mu^{m}\right)_{m=1}^{M}\right\|_{H^{2}(\Omega)} \\
& \quad \leq C_{1}\left(v^{0}\right) C_{2}\left(\left\|u_{I}\right\|_{H^{1}(\Omega)},\left\|\left(B_{V} u_{V}^{m}\right)_{m=1}^{M}\right\|_{L^{2}(\Omega)^{n}},\left\|\left(B_{B} u_{B}^{m}\right)_{m=1}^{M}\right\|_{H^{\frac{1}{2}}(\partial \Omega)^{n}}\right) .
\end{aligned}
$$

Further $e(y, u)$ is Fréchet-differentiable with respect to $y$, and $e_{y}(y, u) \in$ $\mathcal{L}\left(Y_{0}, Z\right)$ has a bounded inverse. Thus Newton's method can be applied for finding the unique solution of 19 for given $u$.

Proof. The existence and stability of the solution for each time instance follows directly from Theorem 4 and Theorem 5

The equation $e(y, u)=0$ is of block diagonal form with nonlinear entries on the diagonal. Thus solving 19 reduces to solving each time instance with given data from the previous time instance. As argued in Theorem 4 and Theorem 5 these nonlinear equations can be solved by Newton's method. Applying this argument for all time instances we obtain that $e_{y}(y, u) \in \mathcal{L}\left(Y_{0}, Z\right)$ has a bounded inverse.

Lemma 7. The functional $e: Y_{0} \times U \rightarrow Z$ is continuously differentiable with respect to $y$ and $u$. Furthermore the equation $e\left(y_{0}, u\right)=0$ for each $u$ admits a unique solution $y(u)$, and $e_{y}\left(y_{0}, u\right)$ is continuously invertible.

The functional $J\left(y_{0}, u\right)$ is continuously differentiable with respect to $y_{0}$ and $u$.

Based on Lemma 7 we introduce the reduced functional $\hat{J}(u):=J\left(y_{0}(u), u\right)$ and state the following theorem.

Theorem 8 (Existence of an optimal control). There exists at least one solution to $\mathcal{P}$, i.e. at least one optimal control.

Proof. Since $\hat{J}$ is bounded from below, there exists a minimizing sequence $u_{l}$ with $\hat{J}\left(u_{l}\right) \rightarrow \hat{J}^{\star}$ and $\hat{J}^{\star}:=\inf _{u} \hat{J}(u)$.

Since $\hat{J}$ is radially unbounded, there exists a closed ball $V \subset U$, bounded, convex and closed such that $u_{l} \subset V$ and thus there exists a weakly convergent subsequence, in the following again denoted by $\left(u_{l}\right)$. Since closed convex sets are weakly closed, $\left(u_{l}\right) \rightharpoonup u_{\star} \in V$ holds. Let $y_{l}=\left(v_{l}, \varphi_{l}, \mu_{l}\right)$ denote the unique solution of (12)-17) for $u_{l}$. Then $y_{l} \rightharpoonup y_{\star} \in Y$, with $y_{\star}=y_{\star}\left(u_{\star}\right)$, and $\left(u_{\star}, y_{\star}\right)$ solves (12)-17. This can be shown as in Garcke et al., 2016. Thm. 6].

It remains to show, that $J\left(u_{\star}\right)=J^{\star}$. Since $y_{l} \rightarrow y_{\star}$ especially $\varphi_{l}^{M} \rightharpoonup \varphi_{\star}^{M}$ and thus by the lower weak semi continuity of norms together with an embedding argument for $W_{u}$ we have

$$
\hat{J}\left(u_{\star}\right) \leq \liminf \hat{J}\left(u_{l}\right)=\hat{J}^{\star} .
$$

Thus $u_{\star}$ is an optimal control.

We next derive first order optimality conditions in the abstract setting. We introduce an adjoint state $p \in Z^{\star}$ and the Lagrangian as

$$
L(y, p, u):=J(y, u)-\langle p, e(y, u)\rangle_{Z^{\star}, Z} .
$$

By Lagrangian calculus we then obtain the following first order optimality conditions.

Theorem 9 (First order optimality conditions in abstract setting). Let $u \in U$, $y \in Y$ be an optimal solution to $\mathcal{P}$. Then there exists an adjoint state $p \in Z^{\star}$ and the triple $(u, y, p)$ fulfills the following first order optimality conditions:

$$
\begin{align*}
e(y, u) & =0 \in Z,  \tag{20}\\
\left(e_{y}(y, u)\right)^{\star} p & =J_{y}(y, u) \in Y_{0}^{\star},  \tag{21}\\
\left\langle J_{u}(y, u)+\left(e_{u}(y, u)\right)^{\star} p, w-u\right\rangle_{U^{\star}, U} & =0 \quad \forall w \in U . \tag{22}
\end{align*}
$$

Proof. From Lemma 7 we have that $e$ and $J$ fulfill the assumptions of Hinze et al., 2009, Cor. 1.3], which in turn asserts the claim.

To state the first order optimality system we introduce Lagrange multiplier $p \in Z^{\star}, p=\left(p_{v}^{m}, p_{\varphi}^{m}, p_{\mu}^{m}\right)_{m=1}^{M} \in H_{0, \sigma}^{M} \times H^{1}(\Omega)^{M} \times H^{1}(\Omega)^{M}$ and define the Lagrangian

$$
\begin{aligned}
L: U & \times\left(H_{0, \sigma}\right)^{M} \times\left(H^{1}(\Omega) \cap L^{\infty}(\Omega)\right)^{M} \times W^{1,3}(\Omega)^{M} \\
& \times\left(H_{0, \sigma}\right)^{M} \times H^{1}(\Omega)^{M} \times H^{1}(\Omega)^{M} \rightarrow \mathbb{R}
\end{aligned}
$$

$$
\begin{aligned}
& L\left(u, v_{0}^{m}, \varphi^{m}, \mu^{m}, p_{v}^{m}, p_{\varphi}^{m}, p_{\mu}^{m}\right):=\frac{1}{2}\left\|\varphi^{M}-\varphi_{d}\right\|_{L^{2}(\Omega)}^{2} \\
& +\frac{\alpha}{2}\left(\alpha_{V}\left\|u_{V}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right)}^{2}+\alpha_{B}\left\|u_{B}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)}^{2}+\alpha_{I}\left(\int_{\Omega} \frac{\delta}{2}\left|\nabla u_{I}\right|^{2}+\frac{1}{\delta} W_{u}\left(u_{I}\right) d x\right)\right) \\
& -\sum_{m=2}^{M}\left[\frac{1}{\tau}\left(\frac{\rho^{m-1}+\rho^{m-2}}{2}\left(v_{0}^{m}+\widetilde{B_{B} u_{B}^{m}}\right)-\rho^{m-2}\left(v_{0}^{m-1}+\widetilde{B_{B} u_{B}^{m-1}}\right), p_{v}^{m}\right)\right. \\
& +a\left(\rho^{m-1}\left(v_{0}^{m-1}+\widetilde{B_{B} u_{B}^{m-1}}\right)+J^{m-1},\left(v_{0}^{m}+\widetilde{B_{B} u_{B}^{m}}\right), p_{v}^{m}\right) \\
& +\left(2 \eta^{m-1} D\left(v_{0}^{m}+\widetilde{B_{B} u_{B}^{m}}\right), D p_{v}^{m}\right) \\
& \left.-\left(\mu^{m} \nabla \varphi^{m-1}, p_{v}^{m}\right)-\left(\rho^{m-1} K, p_{v}^{m}\right)-\left(B_{V} u_{V}^{m}, p_{v}^{m}\right)\right] \\
& -\sum_{m=2}^{M}\left[\frac{1}{\tau}\left(\varphi^{m}-\varphi^{m-1}, p_{\varphi}^{m}\right)+\left(\left(v_{0}^{m}+\widetilde{B_{B} u_{B}^{m}}\right) \nabla \varphi^{m-1}, p_{\varphi}^{m}\right)+\left(b \nabla \mu^{m}, \nabla p_{\varphi}^{m}\right)\right] \\
& -\sum_{m=2}^{M}\left[\sigma \epsilon\left(\nabla \varphi^{m}, \nabla p_{\mu}^{m}\right)-\left(\mu^{m}, p_{\mu}^{m}\right)+\frac{\sigma}{\epsilon}\left(W_{+}^{\prime}\left(\varphi^{m}\right)+W_{-}^{\prime}\left(\varphi^{m-1}\right), p_{\mu}^{m}\right)\right] \\
& -\left[\frac{1}{\tau}\left(\frac{\rho^{1}+\rho^{0}}{2}\left(v^{1}+\widetilde{B_{B} u_{B}^{1}}\right)-\rho^{0} v^{0}, p_{v}^{1}\right)+a\left(\rho^{1} v^{0}+J^{1},\left(v^{1}+\widetilde{B_{B} u_{B}^{1}}\right), p_{v}^{1}\right)\right. \\
& \left.+\left(2 \eta^{1} D\left(v^{1}+\widetilde{B_{B} u_{B}^{1}}\right), D p_{v}^{1}\right)-\left(\mu^{1} \nabla u_{I}, p_{v}^{1}\right)-\left(\rho^{0} K, p_{v}^{1}\right)-\left(B_{V} u_{V}^{1}, p_{v}^{1}\right)\right] \\
& -\left[\frac{1}{\tau}\left(\varphi^{1}-u_{I}, p_{\varphi}^{1}\right)+\left(v^{0} \nabla u_{I}, p_{\varphi}^{1}\right)+\left(b \nabla \mu^{1}, \nabla p_{\varphi}^{1}\right)\right] \\
& -\left[\sigma \epsilon\left(\nabla \varphi^{1}, \nabla p_{\mu}^{1}\right)-\left(\mu^{1}, p_{\mu}^{1}\right)+\frac{\sigma}{\epsilon}\left(W_{+}^{\prime}\left(\varphi^{1}\right)+W_{-}^{\prime}\left(u_{I}\right), p_{\mu}^{1}\right)\right] .
\end{aligned}
$$

Here again $\rho^{m}:=\rho\left(\varphi^{m}\right), \eta^{m}:=\eta\left(\varphi^{m}\right)$ and especially $\rho^{0}:=\rho\left(u_{1}\right), \eta^{0}:=\eta\left(u_{1}\right)$. In the following we write $v^{m}:=v_{0}^{m}+\widehat{B_{B} u_{B}^{m}}$.

The optimality system is now given by $(D L(x), \tilde{x}-x) \geq 0$, where $x$ abbreviates all arguments of $L$ and $\tilde{x}$ denotes an admissible direction. For all components of $x$ except $u_{I}$ it even holds $(D L(x), \tilde{x})=0$ since there no further constraints apply, while $U_{I}$ is a convex subset of $H^{1}(\Omega) \cap L^{\infty}(\Omega)$.

## Derivative with respect to the velocity

The derivative with respect to $v_{0}^{m}$ for $m=2, \ldots, M$ into a direction $\tilde{v} \in H_{0, \sigma}$ is given by

$$
\begin{align*}
\left(D _ { v ^ { m } } L \left(\ldots, v^{m}\right.\right. & , \ldots), \tilde{v})= \\
& \quad-\frac{1}{\tau}\left(\left(\frac{\rho^{m-1}+\rho^{m-2}}{2} \tilde{v}, p_{v}^{m}\right)-\left(\rho^{m-1} \tilde{v}, p_{v}^{m+1}\right)\right) \\
& -a\left(\rho^{m} \tilde{v}, v_{0}^{m+1}+\widehat{B_{B} u_{B}^{m+1}}, p_{v}^{m+1}\right)  \tag{23}\\
& -a\left(\rho^{m-1}\left(v_{0}^{m-1}+\widetilde{B_{B} u_{B}^{m-1}}\right)+J^{m-1}, \tilde{v}, p_{v}^{m}\right) \\
& -\left(2 \eta^{m-1} D \tilde{v}, D p_{v}^{m}\right)-\left(\tilde{v} \nabla \varphi^{m-1}, p_{\varphi}^{m}\right)=0 .
\end{align*}
$$

For $m=1$ we get

$$
\begin{align*}
& \left(D_{v^{1}} L\left(\ldots, v^{1}, \ldots\right), \tilde{v}\right)= \\
& \quad \frac{1}{\tau}\left(\rho^{0} \tilde{v}, p_{v}^{2}\right)-a\left(\rho^{1} \tilde{v}, v_{0}^{2}+\overline{B_{B} u_{B}^{2}}, p_{v}^{2}\right) \\
& \quad-\frac{1}{2 \tau}\left(\left(\rho^{1}+\rho^{0}\right) \tilde{v}, p_{v}^{1}\right)-a\left(\rho^{1}\left(v_{0}^{0}+\widetilde{B_{B} u_{B}^{0}}\right)+J^{1}, \tilde{v}, p_{v}^{1}\right)-\left(2 \eta^{1} D \tilde{v}, D p_{v}^{1}\right)=0 . \tag{24}
\end{align*}
$$

Note that for notational convenience here we introduce artificial variables $v_{0}^{M+1}$, $p_{v}^{M+1}, u_{B}^{M+1}$ and set them to $v_{0}^{M+1} \equiv p_{v}^{M+1} \equiv 0, u_{B}^{M+1}=0$.

Remark 10. Note that we derive the adjoint system in the solenoidal setting. Introducing a variable p for the pressure in the primal equation leads to an additional adjoint variable $p_{p}$ for the adjoint pressure and to an additional term (-div $\left.\tilde{v}, p_{p}\right)$.

## Derivative with respect to the chemical potential

The derivative with respect to the chemical potential for $m=2, \ldots, M$ in a direction $\tilde{\mu} \in W^{1,3}(\Omega)$ is

$$
\begin{align*}
& \left(D_{\mu^{m}} L\left(\ldots, \mu^{m}, \ldots\right), \tilde{\mu}\right)= \\
& \quad-a\left(J_{\mu^{m}}^{m} \tilde{\mu}, v^{m+1}, p_{v}^{m+1}\right)+\left(\tilde{\mu} \nabla \varphi^{m-1}, p_{v}^{m}\right)-\left(b \nabla \tilde{\mu}, \nabla p_{\varphi}^{m}\right)+\left(\tilde{\mu}, p_{\mu}^{m}\right)=0 . \tag{25}
\end{align*}
$$

For $m=1$ the equations is

$$
\begin{align*}
& \left(D_{\mu^{1}} L\left(\ldots, \mu^{1}, \ldots\right), \tilde{\mu}\right)= \\
& \quad-a\left(J_{\mu}^{1} \tilde{\mu}, v^{2}, p_{v}^{2}\right)-a\left(J_{\mu}^{1} \tilde{\mu}, v^{1}, p_{v}^{1}\right)+\left(\tilde{\mu} \nabla u_{I}, p_{v}^{1}\right)-\left(b \nabla \tilde{\mu}, \nabla p_{\varphi}^{1}\right)+\left(\tilde{\mu}, p_{\mu}^{1}\right)=0 . \tag{26}
\end{align*}
$$

Here for $m=1, \ldots, M$ we abbreviate $J_{\mu}^{m} \tilde{\mu}=-\rho_{\delta} b \nabla \tilde{\mu}$, and for notational convenience we introduce artificial variables $v^{M+1}=v_{0}^{M+1}+\widehat{B_{B} u_{B}^{M+1}}, p_{v}^{M+1}$, and set them to $v^{M+1} \equiv p_{v}^{M+1} \equiv 0$.

The above also contains the boundary condition

$$
\nabla p_{\varphi}^{m} \cdot \nu_{\Omega}=0 \quad m=1, \ldots, M
$$

in weak form, which for smooth $p_{\varphi}^{m}$ follows from integration by parts.
Derivative with respect to the phase field
The derivative with respect to the phase field $\varphi^{m}$ in a direction $\tilde{\varphi} \in H^{1}(\Omega) \cap$
$L^{\infty}(\Omega)$ is for $m=2, \ldots, M$

$$
\begin{align*}
& \left(D_{\varphi^{m}} L\left(\ldots, \varphi^{m}, \ldots\right), \tilde{\varphi}\right)= \\
& \quad \delta_{m M}\left(\varphi^{m}-\varphi_{d}, \tilde{\varphi}\right)-\frac{1}{\tau}\left(\rho^{\prime} \frac{v^{m+1} p_{v}^{m+1}+v^{m+2} p_{v}^{m+2}}{2}, \tilde{\varphi}\right)+\frac{1}{\tau}\left(\rho^{\prime} v^{m+1} p_{v}^{m+2}, \tilde{\varphi}\right) \\
& \quad-a\left(\rho^{\prime} \tilde{\varphi} v^{m}, v^{m+1}, p_{v}^{m+1}\right)-\left(2 \eta^{\prime} \tilde{\varphi} D v^{m+1}, D p_{v}^{m+1}\right) \\
& \quad+\left(\mu^{m+1} \nabla \tilde{\varphi}, p_{v}^{m+1}\right)+\left(\rho^{\prime} \tilde{\varphi} K, p_{v}^{m+1}\right) \\
& \quad-\frac{1}{\tau}\left(\left(\tilde{\varphi}, p_{\varphi}^{m}\right)-\left(\tilde{\varphi}, p_{\varphi}^{m+1}\right)\right)-\left(v^{m+1} \nabla \tilde{\varphi}, p_{\varphi}^{m+1}\right) \\
& \quad-\sigma \epsilon\left(\nabla \tilde{\varphi}, \nabla p_{\mu}^{m}\right)-\frac{\sigma}{\epsilon}\left(W_{+}^{\prime \prime}\left(\varphi^{m}\right) \tilde{\varphi}, p_{\mu}^{m}\right)-\frac{\sigma}{\epsilon}\left(W_{-}^{\prime \prime}\left(\varphi^{m}\right) \tilde{\varphi}, p_{\mu}^{m+1}\right)=0 \tag{27}
\end{align*}
$$

where $\delta_{m M}$ denotes the Kronecker delta. For $m=1$ we get

$$
\begin{align*}
\left(D_{\varphi^{1}}\right. & \left.L\left(\ldots, \varphi^{1}, \ldots\right), \tilde{\varphi}\right)= \\
& -\frac{1}{\tau}\left(\frac{\rho^{\prime}}{2} \tilde{\varphi}, v^{2} p_{v}^{2}\right)-a\left(\rho^{\prime} \tilde{\varphi} v^{1}, v^{2}, p_{v}^{2}\right)-a\left(\rho^{\prime} \tilde{\varphi} v^{0}, v^{1}, p_{v}^{1}\right) \\
& -\left(2 \eta^{\prime} \tilde{\varphi} D v^{2}, D p_{v}^{2}\right)-\left(2 \eta^{\prime} \tilde{\varphi} D v^{1}, D p_{v}^{1}\right)-\left(\mu^{2} \nabla \tilde{\varphi} p_{v}^{2}\right)-\left(\rho^{\prime} \tilde{\varphi} K, p_{v}^{2}\right)  \tag{28}\\
& +\frac{1}{\tau}\left(\tilde{\varphi}, p_{\varphi}^{2}\right)-\left(v^{2} \nabla \tilde{\varphi}, p_{\varphi}^{2}\right)-\frac{\sigma}{\epsilon}\left(W_{-}^{\prime \prime}\left(\varphi^{1}\right) \tilde{\varphi}, p_{\mu}^{2}\right) \\
& -\frac{1}{\tau}\left(\frac{\rho^{\prime} \tilde{\varphi}}{2} v^{1}, p_{v}^{1}\right)-\frac{1}{\tau}\left(\tilde{\varphi}, p_{\varphi}^{1}\right)-\sigma \epsilon\left(\nabla \tilde{\varphi}, \nabla p_{\mu}^{1}\right)-\frac{\sigma}{\epsilon}\left(W_{+}^{\prime \prime}\left(\varphi^{1}\right) \tilde{\varphi}, p_{\mu}^{1}\right)=0
\end{align*}
$$

Here for notational convenience we introduce artificial variables $v^{M+1}=v_{0}^{M+1}+$ $\widehat{B_{B} u_{B}^{M+1}}, v^{M+2}=v_{0}^{M+2}+\widehat{B_{B} u_{B}^{M+2}}, p_{v}^{M+1}, p_{v}^{M+2}$, and set them to zero.

The above also contains the boundary condition

$$
\nabla p_{\mu}^{m} \cdot \nu_{\Omega}=0 \quad m=1, \ldots, M
$$

in weak form, which for smooth $p_{\mu}^{m}$ follows from integration by parts.

## Derivative with respect to the control

Finally we calculate the derivative with respect to the control for the three parts of the control space.

For a test direction $w \in U_{V}$ we have

$$
\left(D_{u_{V}} L(u, \ldots), w\right)=\alpha \alpha_{V} \int_{I}\left(u_{V}, w\right)_{\mathbb{R}^{u} v} d t+\sum_{m=1}^{M}\left(B_{V} w^{m}, p_{v}^{m}\right)_{L^{2}(\Omega)}=0
$$

and thus the optimality condition is

$$
\begin{equation*}
\alpha \tau \alpha_{V} u_{V}^{m}+B_{V}^{*} p_{v}^{m}=0 \in \mathbb{R}^{u_{v}} \quad m=1, \ldots, M \tag{29}
\end{equation*}
$$

Here $B_{V}^{\star} p_{v}^{m}$ is defined as

$$
B_{V}^{\star} p_{v}^{m}:=\left(\left(f_{l}, p_{v}^{m}\right)_{L^{2}(\Omega)^{n}}\right)_{l=1}^{u_{v}} .
$$

Concerning the derivative with respect to $u_{B}$ we have for a test function $w \in U_{B}$

$$
\begin{align*}
\left(D_{u_{B}}\right. & L(u, \ldots), w)=\alpha \alpha_{B} \int_{I}\left(u_{B}, w\right)_{\mathbb{R}^{u_{b}}} d t-\tau^{-1}\left(\frac{\rho^{1}+\rho^{0}}{2} \widehat{B_{B} w^{1}}, p_{v}^{1}\right) \\
& -a\left(\rho^{1} v^{0}+J^{1}, \widehat{B_{B} w^{1}}, p_{v}^{1}\right)-2\left(\eta^{1} D \widetilde{B_{B} w^{1}}, D p_{v}^{1}\right) \\
& -\sum_{m=2}^{M}\left[\tau^{-1}\left(\frac{\rho^{m-1}+\rho^{m-2}}{2} \widehat{B_{B} w^{m}}, p_{v}^{m}\right)-\left(\rho^{m-2} \widehat{B_{B} w^{m-1}}, p_{v}^{m}\right)\right.  \tag{30}\\
& +a\left(\rho^{m-1} v^{m-1}+J^{m-1}, \widetilde{B_{B} w^{m}}, p_{v}^{m}\right)+a\left(\rho^{m-1} \widehat{B_{B} w^{m-1}}, v^{m}, p_{v}^{m}\right) \\
& \left.+2\left(\eta^{m-1} D \widetilde{B_{B} w^{m}}, D p_{v}^{m}\right)+\left(\widetilde{B_{B} w^{m}} \nabla \varphi^{m-1}, p_{\varphi}^{m}\right)\right]=0
\end{align*}
$$

For smooth solutions we use the derivative with respect the velocity, the no-flux boundary condition for $v^{m}$ as well as for $\mu^{m}$ and integration by parts to observe

$$
\begin{aligned}
\left(D_{u_{B}} L(u, \ldots), w\right) & =\alpha \alpha_{B} \int_{I}\left(u_{B}, w\right)_{\mathbb{R}^{u_{b}}} d t \\
& -\sum_{m=2}^{M} \int_{\partial \Omega} 2 \eta^{m-1} D p_{v}^{m} \cdot \nu_{\Omega} B_{B} w^{m} d s-\int_{\partial \Omega} 2 \eta^{1} D p_{v}^{1} \cdot \nu_{\Omega} B_{B} w^{1} d s
\end{aligned}
$$

and thus the optimality condition in a strong formulation is

$$
\begin{array}{r}
\alpha \alpha_{B} \tau u_{B}^{m}-\left(\left(2 \eta^{m-1} D p_{v}^{m} \cdot \nu_{\Omega}, g^{l}\right)_{H^{-1 / 2}(\partial \Omega), H^{1 / 2}(\partial \Omega)}\right)_{l=1}^{u_{b}}=0 \in \mathbb{R}^{u_{b}} \quad \forall m=2, \ldots, M, \\
\alpha \alpha_{B} \tau u_{B}^{1}-\left(2 \eta^{1} D p_{v}^{1} \cdot \nu_{\Omega}, g^{l}\right)_{H^{-1 / 2}(\partial \Omega), H^{1 / 2}(\partial \Omega)}=0 \in \mathbb{R}^{u_{b}} . \tag{31}
\end{array}
$$

The derivative with respect to the initial condition $u_{I}$ in a direction $w-u_{I} \in$ $U_{I}$ is

$$
\begin{align*}
\left(D_{u_{I}}\right. & \left.L(u, \ldots), w-u_{I}\right)_{U_{I}^{\star}, U_{I}}=\frac{\alpha}{2} \alpha_{I}\left(\epsilon\left(\nabla u_{I}, \nabla\left(w-u_{I}\right)\right)+\epsilon^{-1} \int_{\Omega} W_{u}^{\prime}\left(u_{I}\right)\left(w-u_{I}\right) d x\right) \\
& -\frac{1}{2 \tau}\left(\rho^{\prime}\left(w-u_{I}\right) v^{2}, p_{v}^{2}\right)+\frac{1}{\tau}\left(\rho^{\prime}\left(w-u_{I}\right) v^{1}, p_{v}^{2}\right) \\
& -\frac{1}{2 \tau}\left(\rho^{\prime}\left(w-u_{I}\right) v^{1}, p_{v}^{1}\right)+\frac{1}{\tau}\left(\rho^{\prime}\left(w-u_{I}\right) v^{0}, p_{v}^{1}\right) \\
& +\left(\mu^{1} \nabla\left(w-u_{I}\right), p_{v}^{1}\right)+\left(\rho^{\prime}\left(w-u_{I}\right) K, p_{v}^{1}\right) \\
& +\frac{1}{\tau}\left(\left(w-u_{I}\right), p_{\varphi}^{1}\right)-\left(v^{0} \nabla\left(w-u_{I}\right), p_{\varphi}^{1}\right)-\frac{\sigma}{\epsilon}\left(W_{-}^{\prime \prime}\left(u_{I}\right)\left(w-u_{I}\right), p_{\mu}^{1}\right) \geq 0 . \tag{32}
\end{align*}
$$

We note that $u_{I} \in H^{1}(\Omega) \cap L^{\infty}(\Omega)$ and thus that there exists no gradient representation for $D_{u_{I}} L$. This is reflected later in our numerical approach.

Remark 11. From (29) we see, that in fact $u_{V}$ has a discrete structure with respect to time, namely it is piecewise constant over time intervals, as the adjoint variable $p_{v}$ is. The same holds for $u_{B}$.

## 4 The fully discrete setting

We next use finite elements to discretize the optimal control problem $\mathcal{P}$ in space. For this we use finite elements on locally adapted meshes. At time instance $t_{m}$, $m=1, \ldots, M$ we use a quasi-uniform, triangulation of $\bar{\Omega}$ with $N T_{m}$ triangles denoted by $\mathcal{T}_{m}=\left\{T_{i}\right\}_{i=1}^{N T_{m}}$ fulfilling $\bar{\Omega}=\bigcup_{i=1}^{N T_{m}} \bar{T}_{i}$.

On $\mathcal{T}_{m}$ we define the following finite element spaces:

$$
\begin{aligned}
& \mathcal{V}_{m}^{1}=\left\{v \in C\left(\mathcal{T}_{m}\right)|v|_{T} \in \mathcal{P}^{1}(T) \forall T \in \mathcal{T}_{m}\right\}, \\
& \mathcal{V}_{m}^{2}=\left\{v \in C\left(\mathcal{T}_{m}\right)^{n}|v|_{T} \in \mathcal{P}^{2}(T)^{n} \forall T \in \mathcal{T}_{m}\right\}
\end{aligned}
$$

where $\mathcal{P}^{l}(S)$ denotes the space of polynomials up to order $l$ defined on $S$. We note that by construction $\mathcal{V}_{m}^{1} \subset W^{1, \infty}\left(\mathcal{T}_{m}\right)$ and $\mathcal{V}_{m}^{2} \subset W^{1, \infty}\left(\mathcal{T}_{m}\right)^{n}$ holds. We introduce the discrete analog to the space $H_{\sigma}(\Omega)$ :

$$
H_{\sigma, m}:=\left\{v \in \mathcal{V}_{m}^{2} \mid(\operatorname{div} v, q)=0 \forall q \in \mathcal{V}_{m}^{1} \cap L^{2}(\Omega)\right\}
$$

and

$$
H_{0, \sigma, m}:=\left\{v \in H_{\sigma, m} \mid \gamma(v)=0\right\} .
$$

We further introduce a linear $H^{1}$-stable projection operator $P^{m}: H^{1}(\Omega) \rightarrow$ $\mathcal{V}_{m}^{1}$ satisfying

$$
\left\|P^{m} v\right\|_{L^{p}(\Omega)} \leq C\|v\|_{L^{p}(\Omega)}, \text { and }\left\|\nabla P^{m} v\right\|_{L^{r}(\Omega)} \leq C\|\nabla v\|_{L^{r}(\Omega)}
$$

for $v \in H^{1}(\Omega)$ with $r \in[1,2]$ and $p \in[1,6]$ if $n=3$, and $p \in[1, \infty)$ if $n=2$ and

$$
\left\|P^{m} v-v\right\|_{H^{1}(\Omega)} \rightarrow 0
$$

for $h \rightarrow 0$ for $v \in H^{2}(\Omega)$ Typically examples are the Clément operator or, by restricting the preimage to $C(\bar{\Omega}) \cap H^{1}(\Omega)$, the Lagrangian interpolation operator.

Using these spaces we state the discrete counterpart of $12-(17)$ :
Let $u \in U$ and $v_{0} \in H_{\sigma} \cap L^{\infty}(\Omega)^{n}$ be given.
Initialization for $m=1$ :
Set $\varphi_{h}^{0}:=u_{I}, v^{0}:=v_{0}$. Find $v_{h}^{1} \in H_{\sigma, 1}, \gamma\left(v_{h}^{1}\right)=\Pi^{1}\left(B_{B} u_{B}^{1}\right), \varphi_{h}^{1} \in \mathcal{V}_{1}^{1}, \mu_{h}^{1} \in \mathcal{V}_{1}^{1}$ such that for all $w \in H_{0, \sigma, 1}, \Psi \in \mathcal{V}_{1}^{1}, \Phi \in \mathcal{V}_{1}^{1}$ it holds:

$$
\begin{align*}
\tau^{-1}\left(\frac{1}{2}\left(\rho_{h}^{1}+\rho_{h}^{0}\right) v_{h}^{1}-\rho_{h}^{0} v^{0}, w\right)+a\left(\rho_{h}^{1} v^{0}+J_{h}^{1}, v_{h}^{1}, w\right) & \\
+\left(2 \eta_{h}^{1} D v_{h}^{1}, D w\right)-\left(\mu_{h}^{1} \nabla \varphi_{h}^{0}+\rho_{h}^{0} g, w\right)-\left(B_{V} u_{V}^{1}, w\right) & =0,  \tag{33}\\
\frac{1}{\tau}\left(\varphi_{h}^{1}-P^{1} \varphi_{h}^{0}, \Psi\right)+\left(b \nabla \mu_{h}^{1}, \nabla \Psi\right)+\left(v^{0} \nabla \varphi_{h}^{0}, \Psi\right) & =0,  \tag{34}\\
\sigma \epsilon\left(\nabla \varphi_{h}^{1}, \nabla \Phi\right)+\frac{\sigma}{\epsilon}\left(W_{+}^{\prime}\left(\varphi_{h}^{1}\right)+W_{-}^{\prime}\left(P^{1} \varphi_{h}^{0}\right), \Phi\right)-\left(\mu_{h}^{1}, \Phi\right) & =0, \tag{35}
\end{align*}
$$

where $J^{1}:=-\rho_{\delta} b \nabla \mu_{h}^{1}$.

Two-step scheme for $m>1$ :
Given $\varphi_{h}^{m-2} \in \mathcal{V}_{m-2}^{1}, \varphi_{h}^{m-1} \in \mathcal{V}_{m-1}^{1}, \mu_{h}^{m-1} \in \mathcal{V}_{m-1}^{1}, v_{h}^{m-1} \in H_{\sigma, m-1}$, find $v_{h}^{m} \in H_{\sigma, m}$, $\gamma\left(v_{h}^{m}\right)=\Pi^{m}\left(B_{B} u_{B}^{m}\right), \varphi_{h}^{m} \in \mathcal{V}_{m}^{1}, \mu_{h}^{m} \in \mathcal{V}_{m}^{1}$ such that for all $w \in H_{0, \sigma, m}, \Psi \in \mathcal{V}_{m}^{1}$, $\Phi \in \mathcal{V}_{m}^{1}$ it holds:

$$
\begin{array}{r}
\tau^{-1}\left(\frac{1}{2}\left(\rho_{h}^{m-1}+\rho_{h}^{m-2}\right) v_{h}^{m}-\rho_{h}^{m-2} v_{h}^{m-1}, w\right)+a\left(\rho_{h}^{m-1} v_{h}^{m-1}+J_{h}^{m-1}, v_{h}^{m}, w\right) \\
+\left(2 \eta_{h}^{m-1} D v_{h}^{m}, D w\right)-\left(\mu_{h}^{m} \nabla \varphi_{h}^{m-1}+\rho_{h}^{m-1} g, w\right)-\left(B_{V} u_{V}^{m}, w\right)=0 \\
\frac{1}{\tau}\left(\varphi_{h}^{m}-P^{m} \varphi_{h}^{m-1}, \Psi\right)+\left(b \nabla \mu_{h}^{m}, \nabla \Psi\right)+\left(v_{h}^{m} \nabla \varphi_{h}^{m-1}, \Psi\right)=0 \\
\sigma \epsilon\left(\nabla \varphi_{h}^{m}, \nabla \Phi\right)+\frac{\sigma}{\epsilon}\left(W_{+}^{\prime}\left(\varphi_{h}^{m}\right)+W_{-}^{\prime}\left(P^{m} \varphi_{h}^{m-1}\right), \Phi\right)-\left(\mu_{h}^{m}, \Phi\right)=0 \tag{38}
\end{array}
$$

where $J_{h}^{m-1}:=-\rho_{\delta} b \nabla \mu_{h}^{m-1}$.
We introduce

$$
\mathcal{V}_{m, b}^{2}:=\left\{\left.v\right|_{\partial \Omega}\left|v \in \mathcal{V}_{m}^{2}, \int_{\partial \Omega} v\right|_{\partial \Omega} \cdot \nu_{\Omega} d s=0\right\}
$$

and define $\Pi^{m}$ for $m=1, \ldots, M$ as the $L^{2}(\partial \Omega)$ projection onto the trace space of $\mathcal{V}_{m, b}^{2}$. This projection is used to incorporate the boundary data and fulfills $\left\|\Pi^{m} g-g\right\|_{L^{2}(\partial \Omega)} \rightarrow 0$ for all $g \in H^{1 / 2}(\partial \Omega)$ with $\int_{\partial \Omega} g \cdot \nu_{\Omega} d s=0$.

We require bounds with respect to $W^{1, p}(\Omega)$-norms for the solution of 33 (38) and prepare these with the following lemmas.

Lemma 12. For all $1<p<\infty$ there exists a continuous function $C(p)$, such that

$$
\|\nabla u\|_{L^{p}(\Omega)} \leq C(p) \sup _{\substack{\left.\eta \in L^{q}(\Omega), \eta \neq 0 \\ \eta, 1\right)=0}} \frac{(\nabla u, \nabla \eta)}{\|\nabla \eta\|_{L^{q}(\Omega)}}
$$

where $\frac{1}{p}+\frac{1}{q}=1$. Further, from the generalized Poincaré inequality, [Alt, 2016, Thm. 8.16], we obtain $\|\eta\|_{W^{1, q}(\Omega)} \leq C\|\nabla \eta\|_{L^{q}(\Omega)}$ and thus

$$
\|\nabla u\|_{L^{p}(\Omega)} \leq C(p) \sup _{\substack{\eta \in L q(\Omega), \eta \neq 0 \\(\eta, 1)=0}} \frac{(\nabla u, \nabla \eta)}{\|\eta\|_{W^{1, q}(\Omega)}}
$$

Proof. The proof follows as in Barrett et al., 2005, Lem. 1.1] and uses $L^{p_{-}}$ stability for $u$ shown in Geng and Shen, 2010, Thm. 1.2].

Lemma 13. For $v \in W^{1, p}(\Omega)$ let $Q_{h} v \in \mathcal{V}_{m}^{1}$ be defined by

$$
\begin{align*}
\left(\nabla Q_{h} v, \nabla w\right) & =(\nabla v, \nabla w) \quad \forall w \in \mathcal{V}_{m}^{1}  \tag{39}\\
\int_{\Omega} Q_{h} v d x & =\int_{\Omega} v d x \tag{40}
\end{align*}
$$

Let $1<p<\infty$. Then it holds

$$
\begin{equation*}
\left\|Q_{h} v\right\|_{W^{1, p}(\Omega)} \leq C(p)\|v\|_{W^{1, p}(\Omega)} . \tag{41}
\end{equation*}
$$

Proof. The proof follows the lines of Brenner and Scott, 2008, Ch. 8]. However, from the fact that the boundary data is of Neumann type new difficulties arise and we refer to Barrett et al., 2005 and Nürnberg and Tucker, 2015 how to deal with these issues.

Lemma 14. Let $u_{h} \in \mathcal{V}_{m}^{1} \subset W^{1, q}(\Omega)$. Then it holds

$$
\left\|\nabla u_{h}\right\|_{L^{p}(\Omega)} \leq C(p) \sup _{\substack{\eta_{h} \in V_{m}^{1}, \eta_{h} \neq 0 \\(\eta, 1)=0}} \frac{\left(\nabla u_{h}, \nabla \eta_{h}\right)}{\left\|\eta_{h}\right\|_{W^{1, q}(\Omega)}}
$$

where $\frac{1}{p}+\frac{1}{q}=1$.
Proof. Directly follows by combining Lemma 12, the definition of $Q_{h} v$ in 40 and the stability estimate (41), compare Nürnberg and Tucker, 2015, Thm. 2.3].

Theorem 15. For given $v_{0} \in H^{1}(\Omega)^{n} \cap L^{\infty}(\Omega)^{n}, u \in U$ there exist $v_{h}^{1} \in H_{\sigma, 1}$, $\gamma\left(v_{h}^{1}\right)=\Pi^{1}\left(B_{B} u_{B}^{1}\right), \varphi_{h}^{1} \in \mathcal{V}_{1}^{1}$ and $\mu_{h}^{1} \in \mathcal{V}_{1}^{1}$ solving (33)-(35). It further holds

$$
\begin{aligned}
& \left\|\mu_{h}^{1}\right\|_{W^{1,3}(\Omega)}+\left\|\varphi_{h}^{1}\right\|_{W^{1,4}(\Omega)}+\left\|v_{h}^{1}\right\|_{H^{1}(\Omega)} \\
& \leq C_{1}\left(v^{0}\right) C_{2}\left(\left\|u_{I}\right\|_{H^{1}(\Omega)},\left\|B_{V} u_{V}^{1}\right\|_{L^{2}(\Omega)^{n}},\left\|B_{B} u_{B}^{1}\right\|_{H^{\frac{1}{2}}(\partial \Omega)^{n}}\right),
\end{aligned}
$$

and Newton's method can be used to find the unique solution to (33-35).
Proof. For (34)-(35) the existence of a unique solution and the applicability of Newton's method follows from (Hintermüller et al., 2011. Also the stability in $H^{1}$ is proven there.

To obtain the estimates of higher regularity we use Lemma 14 . It holds (34)

$$
\begin{align*}
& C\left\|\mu_{h}^{1}\right\|_{W^{1,3}(\Omega)} \leq\left\|\nabla \mu_{h}^{1}\right\|_{L^{3}(\Omega)}+\left\|\mu_{h}^{1}\right\|_{L^{3}(\Omega)} \\
& \quad \leq\left\|\mu_{h}^{1}\right\|_{L^{3}(\Omega)}+C \underset{\substack{v_{h} \in V_{m}^{1},\left(v_{h}, 1\right)=0 \\
\left\|v_{h}\right\|_{W^{1, \frac{3}{2}}(\Omega)}=1}}{ }\left(\nabla \mu_{h}^{1}, \nabla v_{h}\right) \\
& \leq C\left\|\mu_{h}^{1}\right\|_{H^{1}(\Omega)}+C \sup \left(\left|\frac{1}{\tau}\left(\varphi_{h}^{1}-P^{1} \varphi_{h}^{0}, v_{h}\right)\right|+\left|\left(v^{0} \nabla \varphi_{h}^{0}, v_{h}\right)\right|\right) \\
& \leq C\left\|\mu_{h}^{1}\right\|_{H^{1}(\Omega)}+C \sup \left(\left\|\varphi_{h}^{1}-P^{1} \varphi_{h}^{0}\right\|_{L^{2}(\Omega)}\left\|v_{h}\right\|_{L^{2}(\Omega)}+\left\|v^{0} \nabla \varphi_{h}^{0}\right\|_{L^{2}(\Omega)^{n}}\left\|v_{h}\right\|_{L^{2}(\Omega)}\right) \\
& \leq C\left\|\mu_{h}^{1}\right\|_{H^{1}(\Omega)}+C\left\|\varphi_{h}^{1}-P^{1} \varphi_{h}^{0}\right\|_{L^{2}(\Omega)}+C\left\|v^{0} \nabla \varphi_{h}^{0}\right\|_{L^{2}(\Omega)^{n}} \\
& \leq C\left(\left\|\mu_{h}^{1}\right\|_{H^{1}(\Omega)}+\left\|\varphi_{h}^{1}\right\|_{H^{1}(\Omega)}+\left\|\varphi_{h}^{0}\right\|_{H^{1}(\Omega)}+\left\|v^{0}\right\|_{L^{\infty}(\Omega)^{n}}\left\|\varphi_{h}^{0}\right\|_{H^{1}(\Omega)}\right)
\end{align*}
$$

which, together with the already known bound for $\left\|\mu_{h}^{1}\right\|_{H^{1}(\Omega)}$ states the bound on $\mu_{h}^{1}$ in $W^{1,3}(\Omega)$. Note the continuous embedding $W^{1, \frac{3}{2}}(\Omega) \hookrightarrow L^{2}(\Omega)$ used for $v_{h}$.

For $\varphi_{h}^{1}$ we argue similarly and estimate

$$
\begin{aligned}
& C\left\|\varphi_{h}^{1}\right\|_{W^{1,4}(\Omega)} \quad \begin{array}{l}
\leq \varphi_{h}^{1} \|_{L^{4}(\Omega)}+C \sup _{\substack{v_{h} \in W^{1, \frac{4}{3}}(\Omega),\left(v_{h}, 1\right)=0 \\
\| v v_{W^{1, \frac{4}{3}}(\Omega)}}}\left(\left(\nabla \varphi_{h}^{1}, \nabla v_{h}\right)\right) \\
\leq C\left\|\varphi_{h}^{1}\right\|_{H^{1}(\Omega)}+C \sup \left(\left|\left(\mu_{h}^{1}, v_{h}\right)\right|+\left|\frac{\sigma}{\epsilon}\left(W_{+}^{\prime}\left(\varphi_{h}^{1}\right)+W_{-}^{\prime}\left(P^{1} \varphi_{h}^{0}\right), v_{h}\right)\right|\right) \\
\leq C\left\|\varphi_{h}^{1}\right\|_{H^{1}(\Omega)}+C\left\|\mu_{h}^{1}\right\|_{L^{2}(\Omega)}+C \sup \left[\left(1+\left|\varphi_{h}^{1}\right|^{q-1},\left|v_{h}\right|\right)+\left(1+\left|P^{1} \varphi_{h}^{0}\right|^{q-1},\left|v_{h}\right|\right)\right] \\
\leq C\left\|\varphi_{h}^{1}\right\|_{H^{1}(\Omega)}+C\left\|\mu_{h}^{1}\right\|_{L^{2}(\Omega)} \\
\quad+C\left(\left\|1^{2}+\left|\varphi_{h}^{1}\right|^{q-1}\right\|_{L^{2}(\Omega)}+\left\|1+\left|P^{1} \varphi_{h}^{0}\right|^{q-1}\right\|_{L^{2}(\Omega)}\right) \sup \left\|v_{h}\right\|_{L^{2}(\Omega)} \\
\leq C\left\|\varphi_{h}^{1}\right\|_{H^{1}(\Omega)}+C\left\|\mu_{h}^{1}\right\|_{L^{2}(\Omega)}+C\left(1+\left\|\varphi_{h}^{1}\right\|_{H^{1}(\Omega)}+\left\|\varphi_{h}^{0}\right\|_{H^{1}(\Omega)}\right)
\end{array} .
\end{aligned}
$$

We note the continuous embeddings $W^{1, \frac{4}{3}}(\Omega) \hookrightarrow L^{2}(\Omega)$ and $H^{1}(\Omega) \hookrightarrow L^{6}(\Omega)$.
The existence of a unique solution for (33) and stability for $v_{h}^{1}$ then follows from standard arguments for the Oseen equation, since we use an LBB-stable finite element pair.
Theorem 16. For given $u \in U, \varphi^{m-2} \in \mathcal{V}_{m-2}^{1}, \varphi^{m-1} \in \mathcal{V}_{m-1}^{1}, \mu^{m-1} \in \mathcal{V}_{m-1}^{1}$, $v^{m-1} \in H_{\sigma, m-1}$, and for all $m=2, \ldots, M$ there exist $v_{h}^{m} \in H_{\sigma, m}, \gamma\left(v_{h}^{m}\right)=$ $\Pi^{m}\left(B_{B} u_{B}^{m}\right), \varphi_{h}^{m} \in \mathcal{V}_{m}^{1}$ and $\mu_{h}^{m} \in \mathcal{V}_{m}^{1}$ solving (36)-38).

It further holds

$$
\begin{aligned}
& \left\|\mu_{h}^{m}\right\|_{W^{1,3}(\Omega)}+\left\|\varphi_{h}^{m}\right\|_{W^{1,4}(\Omega)}+\left\|v_{h}^{m}\right\|_{H^{1}(\Omega)^{n}} \\
& \leq C\left(\left\|v_{h}^{m-1}\right\|_{H^{1}(\Omega)^{n}},\left\|\mu_{h}^{m-1}\right\|_{W^{1,3}(\Omega)},\left\|\varphi_{h}^{m-1}\right\|_{W^{1,4}(\Omega)}\right. \\
& \left.\left\|B_{V} u_{V}^{m}\right\|_{L^{2}(\Omega)^{n}},\left\|B_{B} u_{B}^{m}\right\|_{H^{\frac{1}{2}}(\partial \Omega)^{n}}\right)
\end{aligned}
$$

and the constant depends polynomially on its arguments.
Proof. In Garcke et al., 2016 the existence of unique solutions to 36 - 38 ) together with bounds in $H^{1}(\Omega)$ on the solution is shown for the case $B_{V} u_{V}^{m}=0$, $B_{B} u_{B}^{m}=0$, using Temam, 1977, Lem. II 1.4]. The volume force $B_{V} u_{V}^{m}$ is given data that enters the proof in a straightforward manner. The boundary data $B_{B} u_{B}^{m}$ can be incorporated by investigating a shifted system as in Theorem 5 .

The estimates of higher regularity follow as in Theorem 15 . There the bound for $\mu_{h}^{1}$ relies on $L^{\infty}(\Omega)$ regularity of $v^{0}$, that is not available here. Instead in (42) we can use a $L^{6}(\Omega)$ bound for $v_{h}^{m}$ that directly follows from the $H^{1}(\Omega)$ bound by Sobolov embedding, together with the $L^{3}(\Omega)$ bound for $\nabla \varphi_{h}^{m-1}$.
Theorem 17. Let $v^{0} \in H^{1}(\Omega)^{n} \cap L^{\infty}(\Omega), u \in U$ be given. Then there exist sequences $\left(v^{m}\right)_{m=1}^{M} \in\left(H_{\sigma, m}\right)_{m=1}^{M},\left(\varphi^{m}\right)_{m=1}^{M},\left(\mu^{m}\right)_{m=1}^{M} \in\left(\mathcal{V}_{m}^{1}\right)_{m=1}^{M}$, such that ( $v^{m}, \varphi^{m}, \mu^{m}$ ) is the unique solution to (33) for $m=1, \ldots, M$. Moreover
there holds

$$
\begin{aligned}
& \left\|\left(v_{h}^{m}\right)_{m=1}^{M}\right\|_{H^{1}(\Omega)}+\left\|\left(\mu_{h}^{m}\right)_{m=1}^{M}\right\|_{W^{1,3}(\Omega)}+\left\|\left(\varphi_{h}^{m}\right)_{m=1}^{M}\right\|_{W^{1,4}(\Omega)} \\
& \leq C_{1}\left(v^{0}\right) C_{2}\left(\left\|u_{I}\right\|_{H^{1}(\Omega)},\left\|\left(B_{V} u_{V}^{m}\right)_{m=1}^{M}\right\|_{L^{2}(\Omega)^{n}},\left\|\left(B_{B} u_{B}^{m}\right)_{m=1}^{M}\right\|_{H^{\frac{1}{2}}(\partial \Omega)^{n}}\right) .
\end{aligned}
$$

Here the constants $C_{1}, C_{2}$ depend polynomially on their arguments.
Proof. The existence of the solution for each time instance follows directly from Theorem 15 and Theorem 16. The stability estimate follows from iteratively applying the stability estimates from Theorem 15 .

Remark 18. The bounds with respect to higher norms are required in Section 4.1 for the limit process $h \rightarrow 0$.

To derive first order necessary optimality conditions we argue as in the case of the time discrete optimization problem and show that Newton's method can be used for solving the primal equation $(33)-\sqrt{38})$ on each time instance.
Theorem 19. Newton's method can be used for finding the unique solution to (33)-(38) on each time instance.

Proof. For $m=1$ this is argued in Theorem 15. For $m>1$ we abbreviate equation (36)-(38) by $F\left(\left(v_{h}^{m}, \varphi_{h}^{m}, \mu_{h}^{m}\right),(w, \Phi, \Psi)\right)=0$. Then $F$ is Fréchet differentiable, since all terms are linear beside the term $W_{+}^{\prime}$ which is differentiable by Assumption (A1) The derivative in the direction $(\delta v, \delta \varphi, \delta \mu) \in H_{0, \sigma, m} \times \mathcal{V}_{m}^{1} \times \mathcal{V}_{m}^{1}$ is given by

$$
\begin{aligned}
& \left\langle G\left(v_{h}^{m}, \varphi_{h}^{m}, \mu_{h}^{m}\right)(\delta v, \delta \varphi, \delta \mu),(w, \Phi, \Psi)\right\rangle:= \\
& \quad \frac{1}{\tau}\left(\frac{\rho^{m-1}+\rho^{m-2}}{2} \delta v, w\right)+a\left(\rho^{m-1} v^{m-1}+J^{m-1}, \delta v, w\right) \\
& \quad+\left(\eta^{m-1} D \delta v, D w\right)-\left(\delta \mu \nabla \varphi^{m-1}, w\right) \\
& \quad+\frac{1}{\tau}(\delta \varphi, \Psi)+(b \nabla \delta \mu, \nabla \Psi)+\left(\delta v \nabla \varphi^{m-1}, \Psi\right) \\
& \quad+\sigma \epsilon(\nabla \delta \varphi, \nabla \Phi)+\frac{\sigma}{\epsilon}\left(W_{+}^{\prime \prime}\left(\varphi_{h}^{m}\right) \delta \varphi, \Phi\right)-(\delta \mu, \Phi) .
\end{aligned}
$$

The existence of a solution ( $\delta v, \delta \varphi, \delta \mu$ ) can be shown following Garcke et al., 2016, Thm. 2], using Brouwer's fixpoint theorem. The boundedness of ( $\delta v, \delta \varphi, \delta \mu$ ) follows from the same proof.

We next introduce the fully discrete analog to problem $(\mathcal{P})$.

$$
\begin{align*}
\min _{u \in U} J\left(\left(\varphi_{h}^{m}\right)_{m=1}^{M}, u\right) & =\frac{1}{2}\left\|\varphi_{h}^{M}-\varphi_{d}\right\|_{L^{2}(\Omega)}^{2} \\
& +\frac{\alpha}{2}\left(\alpha_{I}\left(\int_{\Omega} \frac{\delta}{2}\left|\nabla u_{I}\right|^{2}+\delta^{-1} W_{u}\left(u_{I}\right) d x\right)\right.  \tag{h}\\
& \left.+\alpha_{V}\left\|u_{V}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right)}^{2}+\alpha_{B}\left\|u_{B}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)}^{2}\right)
\end{align*}
$$

s.t. $33-38)$.

We stress, that we do not discretize the control for the initial value. However for a practical implementation we need a discrete description for $u_{I}$. This will be discussed after deriving the optimality conditions, see Section 5

Theorem 20 (Existence of an optimal discrete control). There exists at least one optimal control to $\mathcal{P}_{h}$.

Proof. The claim follows from standard arguments, compare Theorem 8 .

We next state the fully discrete counterpart of the first order optimality conditions from Section 3 .

For this we introduce adjoint variables $\left(p_{v, h}^{m}\right)_{m=1}^{M} \in\left(H_{0, \sigma, m}\right)_{m=1}^{M},\left(p_{\varphi, h}^{m}\right)_{m=1}^{M} \epsilon$ $\left(\mathcal{V}_{m}^{1}\right)_{m=1}^{M}$, and $\left(p_{\mu, h}^{m}\right)_{m=1}^{M} \in\left(\mathcal{V}^{m}\right)_{m=1}^{M}$. For convenience in the following we often write $v_{h}^{m}:=v_{0, h}^{m}+\overline{B_{B} u_{B}^{m}}$.

By the same Lagrangian calculus as in Section 3 we obtain the following fully discrete optimality system.

## Derivative with respect to the velocity

The derivative with respect to $v_{0, h}^{m}$ for $m=2, \ldots, M$ into a direction $\tilde{v} \in \mathcal{V}_{m}^{2}$ is given by

$$
\begin{align*}
& \left(D_{v_{h}^{m}} L\left(\ldots, v_{h}^{m}, \ldots\right), \tilde{v}\right)= \\
& \quad-\frac{1}{\tau}\left(\left(\frac{\rho_{h}^{m-1}+\rho_{h}^{m-2}}{2} \tilde{v}, p_{v, h}^{m}\right)-\left(\rho_{h}^{m-1} \tilde{v}, p_{v, h}^{m+1}\right)\right)  \tag{43}\\
& \\
& \quad-a\left(\rho_{h}^{m} \tilde{v}, v_{h}^{m+1}, p_{v, h}^{m+1}\right)-a\left(\rho_{h}^{m-1} v_{h}^{m-1}+J_{h}^{m-1}, \tilde{v}, p_{v, h}^{m}\right) \\
& \\
& \quad-\left(2 \eta_{h}^{m-1} D \tilde{v}, D p_{v, h}^{m}\right)-\left(\tilde{v} \nabla \varphi_{h}^{m-1}, p_{\varphi, h}^{m}\right)=0 .
\end{align*}
$$

For $m=1$ we get

$$
\begin{align*}
& \left(D_{v_{h}^{1}} L\left(\ldots, v_{h}^{1}, \ldots\right), \tilde{v}\right)= \\
& \quad-\frac{1}{2 \tau}\left(\left(\rho_{h}^{1}+\rho^{0}\right) \tilde{v}, p_{v, h}^{1}\right)+\frac{1}{\tau}\left(\rho^{0} \tilde{v}, p_{v, h}^{2}\right)-a\left(\rho_{h}^{1} \tilde{v}, v_{h}^{2}, p_{v, h}^{2}\right)  \tag{44}\\
& \quad-a\left(\rho_{h}^{1} v_{h}^{0}+J_{h}^{1}, \tilde{v}, p_{v, h}^{1}\right)-\left(2 \eta_{h}^{1} D \tilde{v}, D p_{v, h}^{1}\right)=0 .
\end{align*}
$$

Note that for notational convenience here we introduce artificial variables $v_{h}^{M+1}$, $p_{v, h}^{M+1}$, and set them to $v_{h}^{M+1} \equiv p_{v, h}^{M+1} \equiv 0$.

## Derivative with respect to the chemical potential

The derivative with respect to the chemical potential for $m=2, \ldots, M$ in a direction $\tilde{\mu} \in \mathcal{V}_{m}^{1}$ is

$$
\begin{align*}
& \left(D_{\mu_{h}^{m}} L\left(\ldots, \mu_{h}^{m}, \ldots\right), \tilde{\mu}\right)= \\
& \quad-a\left(J_{\mu} \tilde{\mu}, v_{h}^{m+1}, p_{v, h}^{m+1}\right)+\left(\tilde{\mu} \nabla \varphi_{h}^{m-1}, p_{v, h}^{m}\right)-\left(b \nabla \tilde{\mu}, \nabla p_{\varphi, h}^{m}\right)+\left(\tilde{\mu}, p_{\mu, h}^{m}\right)=0 . \tag{45}
\end{align*}
$$

For $m=1$ the equations is

$$
\begin{align*}
& \left(D_{\mu^{1}} L\left(\ldots, \mu^{1}, \ldots\right), \tilde{\mu}\right)= \\
& \quad-a\left(J_{\mu} \tilde{\mu}, v_{h}^{2}, p_{v, h}^{2}\right)-a\left(J_{\mu}^{1} \tilde{\mu}, v_{h}^{1}, p_{v, h}^{1}\right)+\left(\tilde{\mu} \nabla u_{I}, p_{v, h}^{1}\right)  \tag{46}\\
& \quad-\left(b \nabla \tilde{\mu}, \nabla p_{\varphi, h}^{1}\right)+\left(\tilde{\mu}, p_{\mu, h}^{1}\right)=0
\end{align*}
$$

Here for $m=1, \ldots, M$ we abbreviate $J_{\mu}^{m} \tilde{\mu}=-\rho_{\delta} b \nabla \tilde{\mu}$ and for notational convenience we introduce artificial variables $v_{h}^{M+1}$ and $p_{v, h}^{M+1}$, and set them to $v_{h}^{M+1} \equiv p_{v, h}^{M+1} \equiv 0$.

## Derivative with respect to the phase field

The derivative with respect to the phase field $\varphi_{h}^{m}$ in a direction $\tilde{\varphi} \in \mathcal{V}_{m}^{1}$ is for $m=2, \ldots, M$

$$
\begin{align*}
& \left(D_{\varphi_{h}^{m}} L\left(\ldots, \varphi_{h}^{m}, \ldots\right), \tilde{\varphi}\right)= \\
& \quad \delta_{m M}\left(\varphi_{h}^{m}-\varphi_{d}, \tilde{\varphi}\right)-\frac{1}{\tau}\left(\rho^{\prime} \frac{v_{h}^{m+1} p_{v, h}^{m+1}+v_{h}^{m+2} p_{v, h}^{m+2}}{2}, \tilde{\varphi}\right)+\frac{1}{\tau}\left(\rho^{\prime} v_{h}^{m+1} p_{v, h}^{m+2}, \tilde{\varphi}\right) \\
& \quad-a\left(\rho^{\prime} \tilde{\varphi} v_{h}^{m}, v_{h}^{m+1}, p_{v, h}^{m+1}\right)-\left(2 \eta^{\prime} \tilde{\varphi} D v_{h}^{m+1}, D p_{v, h}^{m+1}\right) \\
& \quad+\left(\mu_{h}^{m+1} \nabla \tilde{\varphi}, p_{v, h}^{m+1}\right)+\left(\rho^{\prime} \tilde{\varphi} g, p_{v, h}^{m+1}\right) \\
& \quad-\frac{1}{\tau}\left(\left(\tilde{\varphi}, p_{\varphi, h}^{m}\right)-\left(P^{m+1} \tilde{\varphi}, p_{\varphi, h}^{m+1}\right)\right)-\left(v_{h}^{m+1} \nabla \tilde{\varphi}, p_{\varphi, h}^{m+1}\right) \\
& \quad-\sigma \epsilon\left(\nabla \tilde{\varphi}, \nabla p_{\mu, h}^{m}\right)-\frac{\sigma}{\epsilon}\left(W_{+}^{\prime \prime}\left(\varphi_{h}^{m}\right) \tilde{\varphi}, p_{\mu, h}^{m}\right)-\frac{\sigma}{\epsilon}\left(W_{-}^{\prime \prime}\left(P^{m+1} \varphi_{h}^{m}\right) P^{m+1} \tilde{\varphi}, p_{\mu, h}^{m+1}\right)=0 \tag{47}
\end{align*}
$$

Here $\delta_{m M}$ denotes the Kronecker delta of $m$ and $M$. For $m=1$ we get

$$
\begin{align*}
\left(D_{\varphi_{h}^{1}}\right. & \left.L\left(\ldots, \varphi_{h}^{1}, \ldots\right), \tilde{\varphi}\right)= \\
& -\frac{1}{\tau}\left(\frac{\rho^{\prime}}{2} \tilde{\varphi}, v_{h}^{2} p_{v, h}^{2}\right)-a\left(\rho^{\prime} \tilde{\varphi} v_{h}^{1}, v_{h}^{2}, p_{v, h}^{2}\right)-a\left(\rho^{\prime} \tilde{\varphi} v^{0}, v_{h}^{1}, p_{v, h}^{1}\right) \\
& -\left(2 \eta^{\prime} \tilde{\varphi} D v_{h}^{2}, D p_{v, h}^{2}\right)-\left(2 \eta^{\prime} \tilde{\varphi} D v_{h}^{1}, D p_{v, h}^{1}\right)-\left(\mu_{h}^{2} \nabla \tilde{\varphi} p_{v, h}^{2}\right)-\left(\rho^{\prime} \tilde{\varphi} g, p_{v, h}^{2}\right) \\
& +\frac{1}{\tau}\left(P^{2} \tilde{\varphi}, p_{\varphi, h}^{2}\right)-\left(v_{h}^{2} \nabla \tilde{\varphi}, p_{\varphi, h}^{2}\right)-\frac{\sigma}{\epsilon}\left(W_{-}^{\prime \prime}\left(P^{2} \varphi_{h}^{1}\right) P^{2} \tilde{\varphi}, p_{\mu, h}^{2}\right) \\
& -\frac{1}{\tau}\left(\frac{\rho^{\prime} \tilde{\varphi}}{2} v_{h}^{1}, p_{v, h}^{1}\right)-\frac{1}{\tau}\left(\tilde{\varphi}, p_{\varphi, h}^{1}\right)-\sigma \epsilon\left(\nabla \tilde{\varphi}, \nabla p_{\mu, h}^{1}\right)-\frac{\sigma}{\epsilon}\left(W_{+}^{\prime \prime}\left(\varphi_{h}^{1}\right) \tilde{\varphi}, p_{\mu, h}^{1}\right)=0 . \tag{48}
\end{align*}
$$

Here for notational convenience we introduce artificial variables $v_{h}^{M+1}, v_{h}^{M+2}$, $p_{v, h}^{M+1}$, and $p_{v, h}^{M+2}$, and set them to zero.
Remark 21. We note that the projection operator $P^{m}$ enters 47) 48 acting on the test function $\tilde{\varphi}$.

## Derivative with respect to the control

Finally we calculate the derivative with respect to the control for the three parts of the control space.

For a test direction $w \in U_{V}$ we have

$$
\left(D_{u_{V}} L(u, \ldots), w\right)=\alpha \alpha_{V} \int_{I}\left(u_{V}, w\right)_{\mathbb{R}^{u_{v}}} d t+\sum_{m=1}^{M}\left(B_{V} w^{m}, p_{v, h}^{m}\right)_{L^{2}(\Omega)}=0
$$

and thus the optimality condition is

$$
\begin{equation*}
\alpha \tau \alpha_{V} u_{V}^{m}+B_{V}^{*} p_{v, h}^{m}=0 \in \mathbb{R}^{u_{v}} \quad m=1, \ldots, M \tag{49}
\end{equation*}
$$

Here $B_{V}^{\star} p_{v}^{m}$ is defined as

$$
B_{V}^{\star} p_{v}^{m}:=\left(\left(f_{l}, p_{v, h}^{m}\right)_{L^{2}(\Omega)^{n}}\right)_{l=1}^{u_{v}} .
$$

Concerning the derivative with respect to $u_{B}$ we have for a test function $w \in U_{B}$

$$
\begin{align*}
\left(D_{u_{B}}\right. & L(u, \ldots), w)=\alpha \alpha_{B} \int_{I}\left(u_{B}, w\right)_{\mathbb{R}^{u_{b}}} d t-\tau^{-1}\left(\frac{\rho_{h}^{1}+\rho^{0}}{2} \widehat{B_{B} w^{1}}, p_{v, h}^{1}\right) \\
& -a\left(\rho_{h}^{1} v^{0}+J_{h}^{1}, \widehat{B_{B} w^{1}}, p_{v, h}^{1}\right)-2\left(\eta_{h}^{1} D \widehat{B_{B} w^{1}}, D p_{v, h}^{1}\right) \\
& -\sum_{m=2}^{M}\left[\tau^{-1}\left(\frac{\rho_{h}^{m-1}+\rho_{h}^{m-2}}{2} \widehat{B_{B} w^{m}}, p_{v, h}^{m}\right)-\left(\rho_{h}^{m-2} \widehat{B_{B} w^{m-1}}, p_{v, h}^{m}\right)\right.  \tag{50}\\
& +a\left(\rho_{h}^{m-1} v_{h}^{m-1}+J_{h}^{m-1}, \widehat{B_{B} w^{m}}, p_{v, h}^{m}\right)+a\left(\rho_{h}^{m-1} \widehat{B_{B} w^{m-1}}, v_{h}^{m}, p_{v, h}^{m}\right) \\
& \left.+2\left(\eta_{h}^{m-1} D \widetilde{B_{B} w^{m}}, D p_{v, h}^{m}\right)+\left(\widetilde{B_{B} w^{m}} \nabla \varphi_{h}^{m-1}, p_{\varphi, h}^{m}\right)\right] \\
& =: \alpha \alpha_{B} \int_{I}\left(u_{B}, w\right)_{\mathbb{R}^{u_{b}}} d t+F_{h}(w)=0
\end{align*}
$$

Here $F_{h}(w)$ abbreviates the action of the discrete normal derivative of $p_{v, h}$, see e.g. Hinze et al., 2009.

The derivative with respect to the initial condition $u_{I}$ in any direction $w-u_{I} \epsilon$ $U_{I}$ is

$$
\begin{align*}
\left(D_{u_{I}}\right. & \left.L(u, \ldots), w-u_{I}\right)_{U_{I}^{\star}, U_{I}}=\frac{\alpha}{2} \alpha_{I}\left(\epsilon\left(\nabla u_{I}, \nabla\left(w-u_{I}\right)\right)+\epsilon^{-1} \int_{\Omega} W_{u}^{\prime}\left(u_{I}\right)\left(w-u_{I}\right) d x\right) \\
& -\frac{1}{2 \tau}\left(\rho^{\prime}\left(w-u_{I}\right) v_{h}^{2}, p_{v, h}^{2}\right)+\frac{1}{\tau}\left(\rho^{\prime}\left(w-u_{I}\right) v_{h}^{1}, p_{v, h}^{2}\right) \\
& -\frac{1}{2 \tau}\left(\rho^{\prime}\left(w-u_{I}\right) v_{h}^{1}, p_{v, h}^{1}\right)+\frac{1}{\tau}\left(\rho^{\prime}\left(w-u_{I}\right) v^{0}, p_{v, h}^{1}\right) \\
& +\left(\mu_{h}^{1} \nabla\left(w-u_{I}\right), p_{v, h}^{1}\right)+\left(\rho^{\prime}\left(w-u_{I}\right) K, p_{v, h}^{1}\right) \\
& \left.+\frac{1}{\tau}\left(\left(w-u_{I}\right), p_{\varphi, h}^{1}\right)-\left(v^{0} \nabla\left(w-u_{I}\right), p_{\varphi, h}^{1}\right)-\frac{\sigma}{\epsilon}\left(W_{-}^{\prime \prime}\left(u_{I}\right)\right)\left(w-u_{I}\right), p_{\mu, h}^{1}\right) \geq 0 \tag{51}
\end{align*}
$$

and this inequality holds for all $w \in U_{I}$.
Remark 22. We use the finite element space $\mathcal{V}_{1}^{1}$ for the representation of $u_{I}$.

### 4.1 The limit $h \rightarrow 0$

We next investigate the limit $h \rightarrow 0$ for problem $\mathcal{P}_{h}$ Let $u^{\star}, \varphi^{\star}$ denote a solution to $\mathcal{P}$ and $u_{h}, \varphi_{h}$ denote a solution to $\mathcal{P}_{h}$. Since $u_{h}, \varphi_{h}$ is a minimizer for $J$ in the discrete setting, we have $J\left(u_{h}, \varphi_{h}\right) \leq J\left(P_{h} u^{\star}, P_{h} \varphi^{\star}\right) \leq C J\left(u^{\star}, \varphi^{\star}\right)=C j$, where $P_{h}$ denotes any $H^{1}$-stable projection onto the discrete spaces. Thus

$$
\begin{align*}
& \frac{1}{2}\left\|\varphi_{h}^{M}-\varphi_{d}\right\|^{2}+\frac{\alpha}{2}\left(\alpha_{I} \int_{\Omega} \frac{\epsilon}{2}\left|\nabla u_{I, h}\right|^{2}+\epsilon^{-1} W_{u}\left(u_{I, h}\right) d x\right.  \tag{52}\\
& \left.\quad+\alpha_{V}\left\|u_{B, h}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right)}^{2}+\alpha_{B}\left\|u_{V, h}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)}^{2}\right) \leq C j
\end{align*}
$$

Note that the mean value of $u_{I, h}$ is fixed and thus by Poincarés inequality we have $\left\|u_{I, h}\right\|_{H^{1}(\Omega)} \leq C\left(1+\left\|\nabla u_{I, h}\right\|\right)$.

Thus from (52) we obtain the following bounds uniform in $h$ :

$$
\left\|u_{I, h}\right\|_{H^{1}(\Omega)}+\left\|u_{B, h}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right)}+\left\|u_{V, h}\right\|_{L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)} \leq C .
$$

Using Theorem 17 we further get the bounds

$$
\left\|\left(v_{h}^{m}\right)_{m=1}^{M}\right\|_{H^{1}(\Omega)^{n}}+\left\|\left(\mu_{h}^{m}\right)_{m=1}^{M}\right\|_{W^{1,3}(\Omega)}+\left\|\left(\varphi_{h}^{m}\right)_{m=1}^{M}\right\|_{W^{1,4}(\Omega)} \leq C .
$$

Using Lax-Milgram's theorem and the above bounds we further obtain bounds

$$
\left\|\left(p_{v, h}^{m}\right)_{m=1}^{M}\right\|_{H^{1}(\Omega)^{n}}+\left\|\left(p_{\varphi, h}^{m}\right)_{m=1}^{M}\right\|_{H^{1}(\Omega)}+\left\|\left(p_{\mu, h}^{m}\right)_{m=1}^{M}\right\|_{H^{1}(\Omega)} \leq C
$$

for the adjoint variables.
Now there exist $u_{I}^{\star} \in H^{1}(\Omega), u_{V}^{\star} \in L^{2}\left(0, T ; \mathbb{R}^{u_{v}}\right), u_{B}^{\star} \in L^{2}\left(0, T ; \mathbb{R}^{u_{b}}\right)$ such that

$$
u_{I, h} \rightharpoonup u_{I}^{\star}, \quad u_{V, h} \rightharpoonup u_{V}^{\star}, \quad u_{B, h} \rightharpoonup u_{B}^{\star} .
$$

There further exist $\left(v^{m, \star}\right)_{m=1}^{M} \in\left(H^{1}(\Omega)^{n}\right)^{M},\left(\varphi^{m, \star}\right)_{m=1}^{M} \in W^{1,4}(\Omega)^{M}$, and $\left(\mu^{m, \star}\right)_{m=1}^{M} \in W^{1,3}(\Omega)^{M}$ such that

$$
v_{h}^{m} \rightharpoonup v^{m, \star}, \quad \varphi_{h}^{m} \rightharpoonup \varphi^{m, \star}, \quad \mu_{h}^{m} \rightharpoonup \mu^{m, \star} \quad \forall m=1, \ldots, M .
$$

And there further exist $\left(p_{v}^{m, \star}\right)_{m=1}^{M} \in H^{1}(\Omega)^{M},\left(p_{\varphi}^{m, \star}\right)_{m=1}^{M} \in H^{1}(\Omega)^{M}$, and $\left(p_{\mu}^{m, \star}\right)_{m=1}^{M} \in$ $H^{1}(\Omega)^{M}$ such that

$$
p_{v, h}^{m} \rightharpoonup p_{v}^{m, \star}, \quad p_{\varphi, h}^{m} \rightharpoonup p_{\varphi}^{m, \star}, \quad p_{\mu, h}^{m} \rightharpoonup p_{\mu}^{m, \star} \quad \forall m=1, \ldots, M
$$

Now let us proceed to the limit in the fully discrete optimality system. To this end we will especially show the following strong convergence results

$$
\begin{array}{cl}
\varphi_{h}^{m} \rightarrow \varphi^{m} & \text { in } H^{1}(\Omega) \\
\mu_{h}^{m} \rightarrow \mu^{m} & \text { in } W^{1,3}(\Omega) \\
v_{h}^{m} \rightarrow v^{m} & \text { in } H_{\sigma}(\Omega) \\
u_{I, h} \rightarrow u_{I} & \text { in } H^{1}(\Omega)
\end{array}
$$

for $m=1, \ldots, M$.

## The limit $h \rightarrow 0$ in the primal equation

The convergence of (38) to (17) and of (35) to (14) follows directly from the proposed weak convergences together with the strong convergence $\varphi_{h}^{m} \rightarrow \varphi^{m}$ in $L^{\infty}$ obtained by compact Sobolev embedding. To obtain strong convergence in $H^{1}(\Omega)$ we argue as in the proof of Theorem 15.

Let $B: H^{1}(\Omega) \times H^{1}(\Omega) \rightarrow \mathbb{R}$ denote the coercive bilinear form $B(u, v)=$ $\sigma \epsilon(\nabla u, \nabla v)+(u, v)$ and let $Q_{h} \varphi^{1} \in \mathcal{V}_{1}^{1}$ denote the projection of $\varphi^{1}$ onto $\mathcal{V}_{1}^{1}$ with respect to $B$ fulfilling $\left\|Q_{h} \varphi^{1}-\varphi^{1}\right\|_{H^{1}(\Omega)} \rightarrow 0$ for $h \rightarrow 0$, since $\varphi^{1} \in H^{2}(\Omega)$.

Then it holds

$$
\left\|\varphi_{h}^{1}-\varphi^{1}\right\|_{H^{1}(\Omega)} \leq\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{H^{1}(\Omega)}+\left\|Q_{h} \varphi^{1}-\varphi^{1}\right\|_{H^{1}(\Omega)}
$$

and

$$
\begin{aligned}
& C\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{H^{1}(\Omega)}^{2} \\
& \quad \leq B\left(\varphi_{h}^{1}-Q_{h} \varphi^{1}, \varphi_{h}^{1}-Q_{h} \varphi^{1}\right)=B\left(\varphi_{h}^{1}-\varphi^{1}, \varphi_{h}^{1}-Q_{h} \varphi^{1}\right) \\
& \leq\left|\left(\mu_{h}^{1}-\mu^{1}, \varphi_{h}^{1}-Q_{h} \varphi^{1}\right)\right|+\left\|\varphi_{h}^{1}-\varphi^{1}\right\|_{L^{2}(\Omega)}\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{L^{2}(\Omega)} \\
& \quad+\frac{\sigma}{\epsilon}\left|\left(W_{+}^{\prime}\left(\varphi_{h}^{1}\right)-W_{+}^{\prime}\left(\varphi^{1}\right), \varphi_{h}^{1}-Q_{h} \varphi^{1}\right)\right|+\frac{\sigma}{\epsilon}\left|\left(W_{-}^{\prime}\left(P^{1} \varphi_{h}^{0}\right)-W_{-}^{\prime}\left(\varphi^{0}\right), \varphi_{h}^{1}-Q_{h} \varphi^{1}\right)\right| \\
& \leq\left\|\mu_{h}^{1}-\mu^{1}\right\|_{L^{2}(\Omega)}\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{L^{2}(\Omega)}+\left\|\varphi_{h}^{1}-\varphi^{1}\right\|_{L^{2}(\Omega)}\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{L^{2}(\Omega)} \\
& \quad+\frac{\sigma}{\epsilon}\left\|W_{+}^{\prime}\left(\varphi_{h}^{1}\right)-W_{+}^{\prime}\left(\varphi^{1}\right)\right\|_{L^{5 / 3}(\Omega)}\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{L^{5 / 2}(\Omega)} \\
& \quad+\frac{\sigma}{\epsilon}\left\|W_{-}^{\prime}\left(P^{1} \varphi_{h}^{0}\right)-W_{-}^{\prime}\left(\varphi^{0}\right)\right\|_{L^{5 / 3}(\Omega)}\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{L^{5 / 2}(\Omega)} .
\end{aligned}
$$

Using Sobolev embedding $H^{1}(\Omega) \hookrightarrow L^{p}(\Omega), p \leq 6$ and dividing by $\left\|\varphi_{h}^{1}-Q_{h} \varphi^{1}\right\|_{H^{1}(\Omega)}$ the resulting differences tend to zero by compact Sobolev embedding, or by Lebesgue's generalized convergence theorem and Assumption (A2). The same arguments apply for the case $m>1$.

The convergence of equation (37) to $\sqrt{16}$ and (34) to $\sqrt{13}$ is shown using the strong convergence $v_{h}^{m} \rightarrow v^{m}$ in $L^{3}(\bar{\Omega})$ together with weak convergence $\nabla \varphi_{h}^{m-1} \rightarrow \nabla \varphi^{m-1}$ in $L^{2}(\Omega)$ yielding weak convergence of the transport term $v_{h}^{m} \nabla \varphi_{h}^{m-1}$ in $L^{6 / 5}$. For $m=1 v^{0} \nabla \varphi_{h}^{0}$ converges weakly in $L^{2}(\Omega)$. Further, strong convergence $\mu_{h}^{m} \rightarrow \mu^{m}$ in $H^{1}(\Omega)$ follows as above.

To show strong convergence in $W^{1,3}$ it is thus sufficient to show strong convergence for $\nabla \mu_{h}^{1} \rightarrow \nabla \mu^{1}$ in $L^{3}(\Omega)$. We define $Q_{h} v \in \mathcal{V}_{1}^{1}$ by

$$
\begin{aligned}
\left(\nabla\left(Q_{h} v-v\right), \nabla w_{h}\right) & =0 \quad \forall w_{h} \in \mathcal{V}_{1}^{1}, \\
\left(Q_{h} v, 1\right) & =(v, 1),
\end{aligned}
$$

satisfying $\left\|Q_{h} v\right\|_{W^{1, \frac{3}{2}}(\Omega)} \leq C\|v\|_{W^{1, \frac{3}{2}}(\Omega)}$, Lemma 13 .

We adapt the idea from Theorem 15 and proceed

$$
\begin{aligned}
& C\left\|\nabla \mu_{h}^{1}-\nabla \mu^{1}\right\|_{L^{3}(\Omega)} \\
& \leq \sup _{\substack{v \in W^{1, \frac{3}{2}}(\Omega),(v, 1)=0 \\
\|v\|_{W^{1,2}}^{2}(\Omega)}}\left(\nabla\left(\mu_{h}^{1}-\mu^{1}\right), \nabla v\right)=\sup \left[\left(\nabla\left(\mu_{h}^{1}-\mu^{1}\right), \nabla Q_{h} v\right)+\left(\nabla\left(\mu_{h}^{1}-\mu^{1}\right), \nabla\left(v-Q_{h} v\right)\right)\right] \\
& \leq \sup \left[\left(\nabla\left(\mu_{h}^{1}-\mu^{1}\right), \nabla Q_{h} v\right)+\left(\nabla \mu^{1}, \nabla\left(Q_{h} v-v\right)\right)\right] \\
& \leq C \sup \left[\left|\tau^{-1}\left(\varphi_{h}^{1}-\varphi^{1}, Q_{h} v\right)\right|+\left|\tau^{-1}\left(P^{1} \varphi_{h}^{0}-\varphi^{0}, Q_{h} v\right)\right|+\mid\left(v^{0} \nabla \varphi_{h}^{0}-v^{0} \nabla \varphi^{0}, Q_{h} v\right)\right. \\
& \left.+\left|\tau^{-1}\left(\varphi^{1}-\varphi^{0}, Q_{h} v-v\right)\right|+\left|\left(v^{0} \nabla \varphi^{0}, Q_{h} v-v\right)\right|\right] \\
& \leq C\left[\left\|\varphi_{h}^{1}-\varphi^{1}\right\|_{L^{2}(\Omega)}+\left\|P^{1} \varphi_{h}^{0}-\varphi^{0}\right\|_{L^{2}(\Omega)}\right. \\
& \left.+\left\|\varphi^{1}-\varphi^{0}\right\|_{L^{2}(\Omega)} \sup \left\|Q_{h} v-v\right\|_{L^{2}(\Omega)}+\left\|v^{0} \nabla \varphi^{0}\right\|_{L^{2}(\Omega)} \sup \left\|Q_{h} v-v\right\|_{L^{2}(\Omega)}\right] \\
& +C \sup \left|\left(v^{0} \nabla Q_{h} v, \varphi_{h}^{0}-\varphi^{0}\right)\right| \\
& \leq C\left[\left\|\varphi_{h}^{1}-\varphi^{1}\right\|_{L^{2}(\Omega)}+\left\|P^{1} \varphi_{h}^{0}-\varphi^{0}\right\|_{L^{2}(\Omega)}\right. \\
& \left.+\left\|\varphi^{1}-\varphi^{0}\right\|_{L^{2}(\Omega)} \sup \left\|Q_{h} v-v\right\|_{L^{2}(\Omega)}+\left\|v^{0} \nabla \varphi^{0}\right\|_{L^{2}(\Omega)} \sup \left\|Q_{h} v-v\right\|_{L^{2}(\Omega)}\right] \\
& +C\left\|v^{0}\right\|_{L^{\infty}(\Omega)}\left\|\varphi_{h}^{0}-\varphi^{0}\right\|_{L^{2}(\Omega)}
\end{aligned}
$$

Note that we used integration by parts to deal with the transport term. From the Hölder and Sobolev inequalities it follows

$$
\left\|Q_{h} v-v\right\|_{L^{2}(\Omega)}^{2} \leq\left\|Q_{h} v-v\right\|_{L^{\frac{3}{2}}(\Omega)}\left\|Q_{h} v-v\right\|_{L^{3}(\Omega)}
$$

The last term is bounded due to the fact, that $\|v\|_{W^{1, \frac{3}{2}}(\Omega)} \leq 1$ and $Q_{h}$ is stable in $W^{1, \frac{3}{2}}(\Omega)$. Since $\left\|Q_{h} v-v\right\|_{L^{\frac{3}{2}}(\Omega)} \leq C h\|v\|_{W^{1, \frac{3}{2}}(\Omega)}$ we obtain $\left\|Q_{h} v-v\right\|_{L^{2}(\Omega)} \rightarrow 0$ for $h \rightarrow 0$ and thus the strong convergence of $\nabla \mu_{h}$ in $L^{3}(\Omega)$. If $m>1$ we can use the strong convergence $\varphi_{h}^{m-1} \rightarrow \varphi^{m-1}$ in $H^{1}(\Omega) \hookrightarrow L^{6}(\Omega)$ together with $\left\|v_{h}^{m}\right\|_{L^{6}(\Omega)} \leq C$ to treat the transport term.

Next we consider the convergence of (36) to (15) and (33) to (12). Here the convergence $\left(\eta_{h}^{m-1} D v_{h}^{m}: D w\right) \rightarrow\left(\eta^{m-1} D v^{m}: D w\right)$ follows from the strong convergence $\varphi_{h}^{m-1} \rightarrow \varphi^{m-1}$ in $L^{\infty}(\Omega)$ (by compact embedding $W^{1,4}(\Omega) \hookrightarrow L^{\infty}(\Omega)$ ) and the weak convergence $D v_{h}^{m} \rightarrow D v^{m}$ in $L^{2}(\Omega)$. The convergence of the trilinear form is obtained by using the just shown strong convergence $\nabla \mu_{h}^{m} \rightarrow \nabla \mu^{m}$ in $L^{3}(\Omega)$ together with the weak convergence of $v_{h}^{m} \rightharpoonup v^{m}$ in $L^{6}(\Omega)$.

Let us finally show strong convergence $v_{h}^{m} \rightarrow v^{m}$ in $H^{1}(\Omega)^{n}$. Let $B: H_{\sigma} \times$ $H_{\sigma} \rightarrow \mathbb{R}$ denote the coercive bilinear form $B(u, v)=\left(\eta_{h}^{1} D u: D v\right)+(u, v)$. The coercivity of $B$ follows from Korn's inequality. Let $w_{h} \in H_{\sigma, 1}$ denote a sequence, such that $\left\|w_{h}-v^{1}\right\|_{H^{1}(\Omega)^{n}} \rightarrow 0$ for $h \rightarrow 0$ and $\left.\left.w_{h}\right|_{\partial \Omega} \equiv v_{h}^{1}\right|_{\partial \Omega}$. The weak continuity of $\Pi^{1}$ ensures that $\left.\Pi^{1}\left(v_{h}^{1}\right) \rightarrow v^{1}\right|_{\partial \Omega}$ and thus such sequence $w_{h}$ exists. Then we have $w_{h} \rightharpoonup v_{h}^{1}$ in $H^{1}(\Omega)$ for $h \rightarrow 0$ and it holds

$$
\left\|v_{h}^{1}-v^{1}\right\|_{H^{1}(\Omega)^{n}} \leq\left\|v_{h}^{1}-w_{h}\right\|_{H^{1}(\Omega)^{n}}+\left\|w_{h}-v^{1}\right\|_{H^{1}(\Omega)^{n}}
$$

Now we proceed with

$$
\begin{aligned}
& C\left\|v_{h}^{1}-w_{h}\right\|_{H^{1}(\Omega)^{n}}^{2} \\
& \quad \leq B\left(v_{h}^{1}-w_{h}, v_{h}^{1}-w_{h}\right)=B\left(v_{h}^{1}, v_{h}^{1}-w_{h}\right)-B\left(w_{h}, v_{h}^{1}-w_{h}\right) \\
& \leq\left(B_{V} u_{h}^{1}, v_{h}^{1}-w_{h}\right)+\left(\mu_{h}^{1} \nabla \varphi_{h}^{0}, v_{h}^{1}-w_{h}\right)+\left(\rho_{h}^{0} g, v_{h}^{1}-w_{h}\right) \\
&-a\left(\rho_{h}^{1} v^{0}+J_{h}^{1}, v_{h}^{1}, v_{h}^{1}-w_{h}\right)-\tau^{-1}\left(\frac{\rho_{h}^{1}+\rho_{h}^{0}}{2} v_{h}^{1}-\rho_{h}^{0} v^{0}, v_{h}^{1}-w_{h}\right) \\
& \quad+\left(v_{h}^{1}, v_{h}^{1}-w_{h}\right)-2\left(\eta_{h}^{1} D w_{h}: D\left(v_{h}^{1}-w_{h}\right)\right)-\left(w_{h}, v_{h}^{1}-w_{h}\right) \\
& \quad \leq\left\|u_{v, h}^{1}\right\|_{L^{2}(\Omega)}\left\|v_{h}^{1}-w_{h}\right\|_{L^{2}(\Omega)}+\left\|\mu_{h}^{1} \nabla \varphi_{h}^{0}\right\|_{L^{\frac{3}{2}}(\Omega)}\left\|v_{h}^{1}-w_{h}\right\|_{L^{3}(\Omega)} \\
& \quad+\left\|\rho_{h}^{0} g\right\|_{L^{2}(\Omega)^{n}\left\|v_{h}^{1}-w_{h}\right\|_{L^{2}(\Omega)}} \quad+\left|a\left(\rho_{h}^{1} v^{0}+J_{h}^{1}, v_{h}^{1}, v_{h}^{1}-w_{h}\right)\right| \\
& \quad+\tau^{-1}\left\|\frac{1}{2}\left(\rho_{h}^{1}+\rho_{h}^{0}\right) v_{h}^{1}-\rho_{h}^{0} v^{0}\right\|_{L^{2}(\Omega)}\left\|v_{h}^{1}-w_{h}\right\|_{L^{2}(\Omega)} \\
& \quad+\left\|v_{h}^{1}-w_{h}\right\|_{L^{2}(\Omega)}^{2}+\left|2\left(\eta_{h}^{1} D w_{h}: D\left(v_{h}^{1}-w_{h}\right)\right)\right| .
\end{aligned}
$$

Now $\left|\left(\eta_{h}^{1} D v^{1}: D\left(v_{h}^{1}-w_{h}\right)\right)\right| \rightarrow 0$ for $h \rightarrow 0$ since $\eta_{h}^{1} D w_{h} \rightarrow \eta^{1} D v^{1}$ in $L^{2}(\Omega)$ and $D\left(v_{h}^{1}-w_{h}\right) \rightharpoonup 0$ in $L^{2}(\Omega)$, and thus beside the trilinear form all terms directly vanish for $h \rightarrow 0$.

For the trilinear form we use the antisymmetry $a\left(\cdot, v_{h}^{1}-w_{h}, v_{h}^{1}-w_{h}\right)=0$ and proceed

$$
\begin{aligned}
\left|a\left(\rho_{h}^{1} v^{0}+J_{h}^{1}, v_{h}^{1}, v_{h}^{1}-w_{h}\right)\right| & =\left|a\left(\rho_{h}^{1} v^{0}+J_{h}^{1}, w_{h}, v_{h}^{1}-w_{h}\right)\right| \\
& =\left|\frac{1}{2}\left(t_{h} \nabla w_{h}, v_{h}^{1}-w_{h}\right)-\frac{1}{2}\left(t_{h} \nabla\left(v_{h}^{1}-w_{h}\right), w_{h}\right)\right| \\
& \leq \frac{1}{2}\left\|t_{h}\right\|_{L^{3}(\Omega)}\left\|w_{h}\right\|_{H^{1}(\Omega)}\left\|v_{h}^{1}-w_{h}\right\|_{L^{2}(\Omega)} \\
& +\frac{1}{2}\left|\left(t_{h} \nabla\left(v_{h}^{1}-w_{h}\right), w_{h}\right)\right|
\end{aligned}
$$

We note the strong convergence $w_{h} \rightarrow v^{1}$ in $L^{6}(\Omega)$ and $t_{h} \rightarrow t$ in $L^{3}(\Omega)$. Thus the last term tends to zero for $h \rightarrow 0$.

For $m>1$ we use $\rho_{h}^{m-1} \rightarrow \rho^{m-1}$ in $L^{\infty}(\Omega)$ to again obtain the strong convergence $\rho_{h}^{m-1} v_{h}^{m-1}+J_{h}^{m-1} \rightarrow \rho^{m-1} v^{m-1}+J^{m-1}$ in $L^{3}(\Omega)$.
The limit $h \rightarrow 0$ in the dual equation
The convergence of (43) and (44) to 23 ) and (24), i.e. the adjoint Navier-Stokes equation, is shown as in the primal equation using the strong convergence of $\varphi_{h}^{m}$ in $L^{\infty}(\Omega)$ and $\mu_{h}^{m}$ in $W^{1,3}(\Omega)$ to show convergence of the trilinear form and of the diffusion term.

The convergence of (45) and (46) to 25 and 26 uses strong convergence of $p_{v, h}^{m+1}$ in $L^{4}(\Omega)$ and of $v_{h}^{m+1}$ in $L^{6}(\Omega)$, where the additional regularity for $v_{h}$ is required.

The convergence of $(47)$ and $(48)$ to $(27)$ and $(28)$ also follows directly using the above shown strong convergence of the primal variables. Especially for the
term $\left(\eta^{\prime} \tilde{\varphi} D v_{h}^{m+1}: D p_{v, h}^{m+1}\right)$ we need the strong convergence $v_{h}^{m+1} \rightarrow v^{m+1}$ in $H^{1}(\Omega)$.
The limit $h \rightarrow 0$ in the derivative w.r.t. the control
The convergence of 49 to 29 is shown using the strong convergence $p_{v, h}^{m}$ in $L^{2}(\Omega)$.

The convergence of (50) to $(30)$ is shown using the various strong convergence results.

Finally we show the convergence of (51) to 32 . Since $J\left(u_{h}, \varphi_{h}\right) \rightarrow J\left(u^{\star}, \varphi^{\star}\right)$, we observe convergence $\left\|\nabla u_{I, h}\right\|_{L^{2}(\Omega)} \rightarrow\left\|\nabla u_{I}^{\star}\right\|_{L^{2}(\Omega)}$. Together with Poincaré's inequality and the weak convergence $u_{I, h} \rightharpoonup u_{I}^{\star}$ in $H^{1}(\Omega)$ we observe strong convergence $u_{I, h} \rightarrow u_{I}^{\star}$ in $H^{1}(\Omega)$. The convergence 51) to 32 now readily follows.

## 5 Numerical examples

In this section we show numerical results for the optimal control problem $\mathcal{P}_{h}$ The implementation is done in C++ using the finite element toolbox FEniCS Logg et al., 2012 together with the PETSc linear algebra backend Balay et al., 2014 and the linear solver MUMPS Amestoy et al., 2001. For the adaptation of the spatial meshes the toolbox ALBERTA [Schmidt and Siebert, 2005] is used. The minimization problem is solved by steepest descent method. If the initial phase field is not used as control, we use the GNU scientific library gsl, 2013, if the initial value is used as control we use a self written implementation using the $H^{1}$ regularity of the control $u_{1}$.

Let us next define some data, that is used throughout all examples. We use $\rho(\varphi)=\frac{\rho_{2}-\rho_{1}}{2} \varphi+\frac{\rho_{1}+\rho_{2}}{2}$ and $\eta(\varphi)=\frac{\eta_{2}-\eta_{1}}{2} \varphi+\frac{\eta_{1}+\eta_{2}}{2}$, where $\rho_{1}, \rho_{2}$ and $\eta_{1}, \eta_{2}$ depend on the actual example. For the free energy we always use (9), with $s=1 e 4$, and the mobility is set to $b \equiv \epsilon / 500$.

### 5.1 The adaptive concept

For the construction of the spatially adapted meshes we use the error indicators that are constructed in Garcke et al., 2016 for the primal equation and use the series of meshes that we construct for the primal equation also for the dual equation. This means that we use classical residual based error estimation to obtain suitable error indicators. We note that following Carstensen and Verfürth, 1999 the cell-wise residuals for the Cahn-Hilliard equation can be subsumed to the edge-wise error indicators. We further note that from our numerical tests we obtain that the cell-wise residuals of the momentum equation is much smaller than the edge-wise indicators, while it turns out to be very expensive to evaluate. Thus we neglect this term. The final error indicator is the cell-wise sum of the jumps of the normal derivatives of the phase field variable, the chemical potential and the velocity field over the cell boundary. The final adaptation scheme for the primal equation is a Dörfler marking scheme based on this indicator, see e.g. Dörfler, 1996 Garcke et al., 2016.

For the Dörfler marking we set the largest cell volume to $V_{\max }=0.0003$, while the smallest cell volume is set to $V_{\min }=\frac{1}{2}\left(\frac{\pi \epsilon}{8}\right)^{2}$ which results in 8 triangles across the interface of thickness $\mathcal{O}(\pi \epsilon)$.

Concerning the temporal resolution, we stress that we did not discretize the control $u_{V}$ and $u_{B}$ with respect to time, i.e. we use the variational discretization approach from Hinze, 2005b. Thus we can adapt the time step size during the optimization to fulfill a CFL-condition without changing the actual control space. Thus we start with a given large time step size $\tau$ and reduce this steps size whenever the CFL-condition $\max _{T} \frac{\left|y^{m}\right|_{T} \mid \tau}{\operatorname{diam(T)}} \leq 1$ is violated for any $m=1, \ldots, M$ by halven $\tau$.

### 5.2 A rising bubble

In this example investigate the pure boundary control $\alpha_{V} \equiv \alpha_{I} \equiv 0$. Here we use $u_{I}=\varphi_{0}$ as given data that we represent on a adapted mesh using the proposed adaptive concept.

We investigate the example of a rising bubble, compare Kahle, 2014 and use the parameters from the benchmark paper Hysing et al., 2009], i.e. $\rho_{1}=$ $1000, \rho_{2}=100, \eta_{1}=10, \eta_{2}=1$. The surface tension is 24.5 which due to our choice of free energy corresponds to $\sigma=15.5972$. The gravitational constant is $g=(0,-0.981)^{t}$ and the computational domain is $\Omega=(0,1) \times(0,1.5)$. The time interval is $I=[0,1.0]$ and we start with a step size $\tau=5 e-3$, that is refined to $\tau=2.5 e-3$ throughout the optimization.

The initial phase field is given by

$$
\varphi_{0}(x)= \begin{cases}\sin \left(\left(\left\|x-M_{1}\right\|-r\right) / \epsilon\right) & \text { if }\left|\left\|x-M_{1}\right\|-r\right| / \epsilon \leq \pi / 2  \tag{53}\\ \operatorname{sign}\left(\left\|x-M_{1}\right\|-r\right) & \text { else }\end{cases}
$$

with $M_{1}=(0.5,0.75)^{t}$ and $r=0.25$. The desired phase field is given by the same expression but with $M_{1}=(0.5,0.5)^{t}$. Thus we aim to move a bubble to the bottom without changing its shape.

Concerning the ansatz functions for the operator $B_{B}$ we introduce the vector field

$$
(f[m, \xi, c](x))_{i}= \begin{cases}\cos \left((\pi / 2)\left\|\xi^{-1}(x-m)\right\|\right)^{2} & \text { if } c \equiv i \text { and }\left\|\xi^{-1}(x-m)\right\| \leq 1 \\ 0 & \text { else }\end{cases}
$$

This describes an approximation to the Gaussian bell with local support. The center is given by $m$ and the diagonal matrix $\xi$ describes the width of the bell in unit directions. We identify a scalar value for $\xi$ with $\xi I$, where $I$ denotes the identity matrix. The parameter $c$ is the number of the component in which the vector field $f$ is non-zero. On the left and right boundary of $\Omega$ we provide 10 equidistantly distributed ansatz functions $f\left[m_{i}, \xi_{i}, c_{i}\right](x)$. Here $\xi_{i}=1.5 / 10$ and $\xi_{i}=1.0 / 10$ if $m_{i}$ is located on bottom or top. We always choose $c_{i}$ such that the ansatz function is tangential to $\Omega$.


Figure 1: The initial phase field $\varphi_{0}$ (left), the desired phase field $\varphi_{d}$ (middle) and the control areas together with the zero-level lines of $\varphi_{0}$ and $\varphi_{d}$ (right) for the rising bubble example. Note that each of the control areas contains 10 controls of the type $f[m, \xi, c](x)$ that point tangential to $\partial \Omega$ with overlapping support.

We set $\alpha=1 e-10$ and $\epsilon=0.04$ and stop the optimization as soon as $\|\nabla J(u)\|_{U}$ is decreased by a factor of 0.1 .

In Figure 1 we present the initial phase field $\varphi_{0}$, the desired phse field $\varphi_{d}$ and the control areas together with the zero-level lines of $\varphi_{0}$ and $\varphi_{d}$.

The steepest descent method is able to reduce $\|\nabla J\|_{U}$ from $6 e-2$ to $4.6 e-2$ in 67 iterations and stagnates due to no further decrease in $\|\nabla J\|_{U}$. Mean while the functional $J$ is reduced from 0.509 to 0.033 . In Figure 2 we show the evolution of $\varphi$ for the optimal control together with the magnitude of the velocity field.

In Figure 3 we show the evolution of the control action over time. We observe a rapid decay of the control strength at the end of the time horizon, while the first peak corresponds to a strong control at the side walls in the region above the bubble, that is rather inactive after this initial stage.

### 5.3 Reconstruction of the initial value

Finally we investigate an example of finding an initial phase field, such that after a given amount of time without further control action a desired phase field is achieved. Here we apply only initial value control, i.e. $\alpha_{V}=\alpha_{B}=0$, and we use no-slip boundary conditions for the velocity field.

Let us turn to the representation of $u_{I}$. We initialize $u_{I}$ with a constant value $u_{I}=-0.8$ and use a homogeneously refined initial mesh for its representation. We use this mesh for $\mathcal{T}_{1}$.

After each step of the minimization algorithm we use the jumps accross edges in normal direction of $\nabla u_{I}$ to construct a new grid for the representation of $u_{I}$ and interpolate the current control to the new grid. The marking is evaluated based on a Dörfler approach.

The parameter for this example are given as $\rho_{1}=1000, \rho_{2}=1, \eta_{1}=10$,


Figure 2: The evolution of the optimally controlled phase field and velocity field at times $t=0.25,0.5,0.75,1.0$ (left to right) when control is only applied to the side walls and not at the bottom and the top part of the boundary. The pictures show the magnitude of the velocity field on the left and the phase field on the right. For $t=0.1$ we additionally indicate the zero-level line of $\varphi_{d}$ by a black line. Note that the velocity field coincides with $B_{B} u_{B}$ on the boundary.


Figure 3: The evolution of the optimal control action over time, i.e. $\|u(t)\|$ for the rising bubble example.


Figure 4: The optimal control $u_{I}$ (left) and the resulting distribution at the end of the time interval (right), where the bubble is shown in gray. The black line indicates the zero level line of the desired shape.
$\eta_{2}=0.1, \sigma=1.245$ and $g \equiv-0.981$. These are the parameters of the second benchmark from Hysing et al., 2009, where $\sigma$ was rescaled due to our specific choice of energy. We note that due to the large ratio in density, the bubble undergoes strong deformation during rising. The optimization horizon again is $I=[0,1.5]$, and $\Omega=(0,1)^{2}$. We set $\alpha=0.2$ and solve the optimization problem for $\epsilon=0.02$.

We initialize the optimization with $u_{I} \equiv-0.8$ and use a circle around $M=$ $(0.5,0.6)$ with radius $r=0.1763040551$ as defined in 53 ) as desired shape. These values are used such that $\int_{\Omega} \varphi_{d}-u_{I} d x=0$ is fulfilled.

The optimization problem is solved using the VMPT method, proposed in Blank and Rupprecht, 2015. It is an extension of the projected gradient method to the Banach space setting. In our situation this is $H^{1}(\Omega) \cap L^{\infty}(\Omega)$.

We stop the allover algorithm as soon as $\left|\left(D J_{u_{I}}(\cdot), v\right)\right|<1 e-3$, where $v$ denotes the current normalized search direction. In our example this is reached after 31 iterations, where $J$ is reduced from $3.8 \mathrm{e}-1$ to $1.9 \mathrm{e}-1$, and especially $\left\|\varphi^{K}-\varphi_{d}\right\|$ is reduced from 0.43 to 0.16 .

In Figure 4 we show the initial shape at the end of the optimization process, on the left and the corresponding shape at the end of the optimization time interval together with the zero level line of the desired shape on the right.

Remark 23. In first examples we used an energy for $W_{u}$ that fulfills Assumptions (A1) (A4) and the method of steepest descent to solve the resulting optimization problem. There we only got very slow convergence of the algorithm and the resulting optimal $u_{I}$ had much broader interfaces. So it seems that it is recommended to use the non-smoth free energy as we propose here.

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