# A Smooth Vector Field for Saddle Point Problems 

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#### Abstract

In this paper we propose a novel smooth vector field whose trajectories globally converge to the saddle point of the Lagrangian associated with a convex and constrained optimization problem. Under suitable assumptions, we prove global convergence of the trajectories for the class of strictly convex problems and we propose a vector field for linear programs.


## I. INTRODUCTION

In this paper we propose a novel way to compute saddle points that arise in convex optimization problems of the form

$$
\begin{align*}
& \inf _{x} f(x)  \tag{1}\\
& \text { s.t. } g_{i}(x) \leq 0, i=1, \ldots, m
\end{align*}
$$

with $x \in \mathbb{R}^{n}, f: \mathbb{R}^{n} \rightarrow \mathbb{R}, g_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}, f, g_{i}$ convex. Under the assumption of strict feasibility (Slater's Condition), i.e. there exists an $\tilde{x}$ with $g_{i}(\tilde{x})<0$ for $i=1, \ldots, m, x^{*}$ is the solution of (1) if and only if $\left(x^{*}, \lambda^{*}\right)$ is a saddle point of the Lagrangian, i.e. $L\left(x^{*}, \lambda\right) \leq L\left(x^{*}, \lambda^{*}\right) \leq L\left(x, \lambda^{*}\right)$ for all $x \in \mathbb{R}^{n}, \lambda \in \mathbb{R}_{+}^{m}$ with

$$
\begin{equation*}
L(x, \lambda)=f(x)+\sum_{i=1}^{m} \lambda_{i} g_{i}(x) \tag{2}
\end{equation*}
$$

We propose a continuous-time dynamical system to find a saddle point of the Lagrangian (2). The idea of using a dynamical system to find a saddle point of a Lagrangian function goes back to K. J. Arrow, L. Hurwicz and H. Uzawa as well as to G. W. Brown and J. von Neumann who proposed in [2] (see also [17]) and [7] respectively, a gradient-like system. In order to assure the satisfaction of the constraints and the positivity of the Lagrange multipliers, a method involving a projected gradient in the vector field, i.e.

$$
\begin{align*}
\dot{x} & =-\left(\frac{\partial L(x, \lambda)}{\partial x}\right)^{\top}=-\nabla f(x)-\sum_{i=1}^{m} \lambda_{i} \nabla g_{i}(x)  \tag{3a}\\
\dot{\lambda}_{i} & =\mathcal{P}\left(\lambda_{i}, g_{i}(x)\right), i=1, \ldots, m \tag{3b}
\end{align*}
$$

where the operator $\mathcal{P}\left(\lambda_{i}, g_{i}(x)\right):=0$ if $\lambda_{i}=0$ and $g_{i}(x)<0$, and $\mathcal{P}\left(\lambda_{i}, g_{i}(x)\right):=g_{i}(x)$ otherwise, was used to project the vector field to the cone of positive vectors in $\mathbb{R}_{+}^{m}$. The projection as it was done in the Arrow-Hurwicz-Uzawa flow (AHU-flow) in (3) renders the vector field non-smooth. Since the work of Arrow, Hurwicz and Uzawa, a lot of further research in the direction of saddle point algorithms has been done, especially in the areas of economics, optimization

[^0]and game theory. For example, the so-called best response dynamics [4], [14] is a continuous-time dynamical system of the form $\dot{x} \in B R_{1}(y)-x, \dot{y} \in B R_{2}(x)-y$ which converges to a saddle point of a zero sum game. While this system is of interest for several reasons, it does not directly yield to an easy implementable algorithm for finding saddle points, since the evaluation of the best response $B R_{1}(\cdot)$ is an optimization problem itself.

Increasing interest in distributed control and optimization has recently motivated more research in the direction of saddle point algorithms (see [10] and references therein). For example, in [20] a discrete-time primal-dual subgradient algorithm for saddle point problems has been proposed which builds upon the AHU-flow. As pointed out in [20], primal-dual methods are especially appealing in distributed optimization over networks when the dual function cannot be evaluated efficently. The algorithm proposed in this work is also a primal-dual method, but it is in the following aspects different from the work outlined above.

Firstly, in contrast to the AHU-flow, we propose a new method to assure positivity of the Lagrange multipliers by avoiding the (non-smooth) projection operator. Moreover, our method can easily be extended to other cases, e.g. when the Lagrange multipliers are elements in the positive semidefinite matrix cone. Furthermore, the overall idea behind the proposed dynamics is different from the AHU-flow and implicitly involves the dual function. Concerning convergence, we establish global convergence to a saddle point under comparable assumptions as in the AHU-flow. Additionally, the proposed dynamics shows promising results for linear programs and guarantees global asymptotic stability in that case.

Secondly, in contrast to the best response dynamics, the evaluation of the right hand-side involves no optimization problem and is thus easy to implement. What is needed, however, is the evaluation of a Hessian, which is dropped in the linear programming case.

Finally, in contrast to the work of Nedić and Ozdaglar (see [20]), our algorithm is formulated in continuous-time. While discrete-time formulations are as powerful as a continuoustime formulation and often more convenient to implement on a computer, it is on the other hand often advantageous to work in continuous-time, since many powerful mathematical concepts and tools can be applied. Moreover, a differential equation is independent of its algorithmic realization, whereas an algorithm given in the form $x_{k+1}=x_{k}+\alpha f\left(x_{k}\right)$ has a fixed realization.

Summarizing, our main contribution in this paper is a dynamical system, given as differential equation, which
complements the more dominating literature on discrete-time saddle point algorithms and provides a novel, alternative approach to the AHU-flow with guaranteed convergence properties by avoiding the use of projection operators in the vector field. We prove global convergence of the trajectories to a - not necessarily unique - saddle point for a certain class of constrained optimization problems and propose a modified version for linear programs.

The remainder of this paper is structured as follows. In Section II we introduce the notation and the necessary mathematical preliminaries. In Section III we state our main result and the corresponding proofs. In Section IV we demonstrate some aspects of the algorithm by examples. Finally, in Section V we conclude with a summary and outlook.

## II. BACKGROUND

We will make use of the following notation:
The norm of a vector $v \in \mathbb{R}^{n}$ will be denoted by $\|v\|=$ $\sqrt[2]{\sum_{i=1}^{n} v_{i}^{2}}$. A matrix $A$ is said to be positive semidefinite (definite), or $A \geq 0(A>0)$, if it is symmetric $A=A^{\top}$, and all its eigenvalues are nonnegative (positive). A matrix $H$ is Hurwitz if the real parts of its eigenvalues are in the complex open left half-plane.
$\mathbb{R}_{+}^{m}$ denotes the cone of nonnegative vectors in $\mathbb{R}^{m}$ and $\mathbb{R}_{++}^{m}$ the interior of $\mathbb{R}_{+}^{m}$, i.e. the set of positive vectors in $\mathbb{R}^{m}$. Define the operator diag : $\mathbb{R}^{n} \rightarrow \mathbb{R}^{n \times n}$ that takes a vector and maps it to a matrix with the elements of the vector on its main diagonal and define the open ball $B_{\epsilon}\left(x_{0}\right)$ as the set $B_{\epsilon}\left(x_{0}\right)=\left\{x \in \mathbb{R}^{n}:\left\|x-x_{0}\right\|<\epsilon\right\}$.

Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ and $L: \mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}$ be a $C^{1}$ - and a $C^{2}$-function, respectively. Then

$$
\begin{aligned}
\nabla f(x) & =\left(\frac{\partial f(x)}{\partial x}\right)^{\top}=\left[\frac{\partial f(x)}{\partial x_{1}}, \ldots, \frac{\partial f(x)}{\partial x_{n}}\right]^{\top} \\
\nabla L(x, \lambda) & =\left(\frac{\partial L(x, \lambda)}{\partial x}\right)^{\top}=\left[\frac{\partial L(x, \lambda)}{\partial x_{1}}, \ldots, \frac{\partial L(x, \lambda)}{\partial x_{n}}\right]^{\top} \\
\nabla^{2} L(x, \lambda) & =\frac{\partial^{2} L(x, \lambda)}{\partial x^{2}}
\end{aligned}
$$

We continue with some standard results from convex optimization. The following three statements are equivalent: $f$ convex $\Leftrightarrow \nabla f(x)^{\top}(y-x) \leq f(y)-f(x) \Leftrightarrow \nabla^{2} f(x) \geq 0$. Furthermore we have $\forall x \neq y$ : $f$ strictly convex $\Leftrightarrow$ $\nabla f(x)^{\top}(y-x)<f(y)-f(x)$ and on the other hand, a positive definite Hessian $\nabla^{2} f(x)>0$ implies strict convexity of $f$.

A saddle point of (2) delivers always a solution to (1). In order to guarantee that the converse is true, it suffices to guarantee strict feasibility (Slater's Condition).

The following theorems can be found in the standard literature dealing with optimization. See for example [5], [9] and [21].

Theorem 1: Suppose Slater's Condition is satisfied and $f, g_{i}$ for $i=1, \ldots, m$ are convex. Then $x^{*}$ is a solution of (1) if and only if there exists a $\lambda^{*}$ such that $\left(x^{*}, \lambda^{*}\right)$ is a saddle point of (I) in $\mathbb{R}^{n} \times \mathbb{R}_{+}^{m}$, i.e. $\left(x^{*}, \lambda^{*}\right)=$ $\arg \inf _{x \in \mathbb{R}^{n}} \sup _{\lambda \in \mathbb{R}_{+}^{m}} L(x, \lambda)$.

Another closely related result is delivered by the Kuhn-Karush-Tucker Theorem.

Theorem 2: Suppose Slater's Condition is satisfied and $f, g_{i} \in C^{1}$ for $i=1, \ldots, m$ are convex. Then $x^{*}$ is a solution of (1) if and only if there exists a $\lambda^{*} \in \mathbb{R}_{+}^{m}$ such that the following three conditions are satisfied:

$$
\begin{align*}
& \nabla f\left(x^{*}\right)+\sum_{i=1}^{m} \lambda_{i}^{*} \nabla g_{i}\left(x^{*}\right)=0  \tag{4a}\\
& \lambda_{i}^{*} g_{i}\left(x^{*}\right)=0, i=1, \ldots, m  \tag{4b}\\
& g_{i}\left(x^{*}\right) \leq 0, i=1, \ldots, m \tag{4c}
\end{align*}
$$

The conditions are also often referred to as KKT-Conditions.

## III. Main Results

We propose the following equations to find a saddle point of the Lagrangian (2):

$$
\begin{align*}
\dot{x} & =-\nabla L(x, \lambda)-\nabla^{2} L(x, \lambda)^{-1} \sum_{i=1}^{m} \lambda_{i} g_{i}(x) \nabla g_{i}(x)  \tag{5a}\\
\dot{\lambda}_{i} & =\lambda_{i} g_{i}(x), i=1, \ldots, m \tag{5b}
\end{align*}
$$

with initial conditions $x(0)=x_{0} \in \mathbb{R}^{n}$ and $\lambda_{i}(0)=\lambda_{i 0} \in$ $\mathbb{R}_{++}$for $i=1, \ldots, m$.

Before we prove the main results, we briefly discuss the structure of these equations and the underlying ideas.

Firstly, we observe that the equilibrium points of (5a) and (5b) in $\mathbb{R}^{n} \times \mathbb{R}_{+}^{m}$ are the points which satisfy the KKTConditions (4a) and (4b).

Secondly, comparing the equations with the AHU-flow (3), we see a modification in the $\lambda$-dynamics in (5b). The idea behind (5b) is to guarantee non-negativity of the Lagrange multipliers in a natural way and to avoid switching (projection) as it was done in the AHU-flow. Notice that the $\lambda_{i}$ 's stay non-negative if initialized in $\mathbb{R}_{++}^{m}$. It is also not difficult to see that (5b) can be generalized to other conic constraints, e.g. when dealing with positive semidefinite Lagrange multiplier matrices. In that case one can imagine to replace (5b) by the matrix equation $\dot{\Lambda}=\frac{1}{2}\left(G^{\top} \Lambda+\Lambda G\right)$ with $\Lambda>0$ (for further information on dynamical systems involving matrices see e.g. [1] and [13]).

Thirdly, we see an additional expression on the right hand side in the $x$-dynamics (5a). This expression arises from the following idea which leads to the system in (5). The task of the $x$-dynamics (5a) is to minimize $L(x, \lambda)$. If $\lambda(t)$ was constant, then the right hand side of (5a) would just consist of the negative gradient of $L(x, \lambda)$, but since $\lambda(t)$ itself is evolving, its dynamics have to be taken into account, which leads to the additional term $-\nabla^{2} L(x, \lambda)^{-1} \sum_{i=1}^{m} \lambda_{i} g_{i}(x) \nabla g_{i}(x)$.
The overall idea underlying the dynamics of system (5) is illustrated in Fig. 1. The system is designed in such a way that the $x$-dynamics will minimize $L(x, \lambda)$ with respect to $x$, i.e. by zeroing

$$
\begin{equation*}
V_{1}(x, \lambda)=\frac{1}{2} \nabla L(x, \lambda)^{\top} \nabla L(x, \lambda) \tag{6}
\end{equation*}
$$

while the $\lambda$-dynamics will maximize the (dual) function

$$
\begin{equation*}
V_{2}(\lambda)=-L(\bar{x}(\lambda), \lambda) \tag{7}
\end{equation*}
$$



Fig. 1: Tracking of $\bar{x}(\lambda)$
on $\mathbb{R}_{+}^{m}$, with $\bar{x}(\lambda)$ as the minimizer $x$ of $L(x, \lambda)$ with respect to some $\lambda \in \mathbb{R}_{+}^{m}$. In the following, we prove convergence to a saddle point of $L(x, \lambda)$ in (2), assuming that Slater's Condition is satisfied, $f, g_{i}$ 's are strictly convex and possess a minimum.

Hereby, we proceed in several steps. In Lemma 1 we show that $\bar{x}_{c}(\lambda)$ (see (8), (9), (10)) with respect to $\lambda \in \mathbb{R}_{+}^{m}$ is always bounded. In Lemma 2 we establish radially unboundedness of $V_{1}(x, \lambda)$ with respect to $x$ by introducing an auxiliary optimization problem. Lemma 3 states that under the assumption of forward completeness of $(x(t), \lambda(t)), x(t)$ converges to the set $V_{1}(x, \lambda)=0$ whereas the existence and boundedness of $x(t)$ and $\lambda(t)$ is established in Lemma 4. Finally in Theorem 3 we prove convergence of trajectories to a saddle point.

Global convergence of the trajectories to a saddle-point is established under the assumptions (A1) - (A3).

## Assumptions:

(A1) $f, g_{i}, i=1, \ldots, m$, are twice continuously differentiable and strictly convex. Moreover $\nabla^{2} f(x)>0$.
(A2) The functions $f_{c}(x)=f(x)-c^{\top} x, g_{i}, i=1, \ldots, m$ possess a minimum for any $c \in \mathbb{R}^{n}$.
(A3) $\exists \tilde{x} \in \mathbb{R}^{n}$ such that $g_{i}(\tilde{x})<0, i=1, \ldots, m$.
Assumption (A2) is necessary since strict convexity of a function does not imply compact level sets, as e.g. the function $-\ln (x)$ is strictly convex, but does not admit a minimum. Since we need compact level sets, we restrict $f$ and all $g_{i}$ 's to be strictly convex and additionally $f_{c}$ and all $g_{i}$ 's must admit a minimum. The additional linear term $-c^{\top} x$ appearing in $f_{c}$ is a technical assumption needed in the proof and assures, loosely speaking, that $f$ grows faster than a linear function.

We now introduce an auxiliary optimization problem, whose link to problem (1) will be clear immediately. Consider the following optimization problem

$$
\begin{align*}
& \inf _{x} f(x)-c^{\top} x  \tag{8}\\
& \text { s.t. } g_{i}(x) \leq 0, i=1, \ldots, m
\end{align*}
$$

with $c \in \mathbb{R}^{n}$. Observe that the Lagrangian

$$
\begin{equation*}
L_{c}(x, \lambda)=f(x)-c^{\top} x+\sum_{i=1}^{m} \lambda_{i} g_{i}(x) \tag{9}
\end{equation*}
$$

associated with problem (8) is strictly convex in $x$. Before we state our main results, we introduce the quantity $\bar{x}_{c}(\lambda)$
that is implicitly defined by
$\nabla L_{c}\left(\bar{x}_{c}(\lambda), \lambda\right)=\nabla f\left(\bar{x}_{c}(\lambda)\right)-c+\sum_{i=1}^{m} \lambda_{i} \nabla g_{i}\left(\bar{x}_{c}(\lambda)\right)=0$.
Obviously, the Lagrangian $L(x, \lambda)$ in (2) associated with problem (1) coincides with the Lagrangian $L_{c}(x, \lambda)$ in (9) for $c=0$, i.e. $L(x, \lambda)=L_{0}(x, \lambda)$. Therefore, we can also state that $\bar{x}(\lambda)$ as the minimum of $L(x, \lambda)$ for some $\lambda \in \mathbb{R}_{+}^{m}$ coincides with $\bar{x}_{0}(\lambda)$, i.e. $\bar{x}(\lambda)=\bar{x}_{0}(\lambda)$.
$L_{c}(x, \lambda)$ is strictly convex in $x$. This implies that $\bar{x}_{c}(\lambda)$ is the unique solution of the implicit equation (10) for any $\lambda \in \mathbb{R}_{+}^{m}$ and $c \in \mathbb{R}^{n}$. Due to assumption (A2) and the Implicit Function Theorem, $\bar{x}_{c}(\lambda)$ is finite for any finite $\lambda \in \mathbb{R}_{+}^{m}$ and is a $C^{1}$-function. Its derivative is given by (see Appendix VI-A) $\frac{\partial \bar{x}_{c}(\lambda)}{\partial \lambda}=$ $-\left(\nabla^{2} L\left(\bar{x}_{c}(\lambda), \lambda\right)\right)^{-1}\left[\nabla g_{1}\left(\bar{x}_{c}(\lambda)\right), \ldots, \nabla g_{m}\left(\bar{x}_{c}(\lambda)\right)\right]$.

Lemma 1: Let assumptions (A1) - (A3) be satisfied and suppose $\bar{x}_{c}(\lambda)$ is the unique minimizer of $L_{c}(x, \lambda)$ for a given $\lambda \in \mathbb{R}_{+}^{m}$ and $c \in \mathbb{R}^{n}$. Then $\sum_{i=1}^{m} \lambda_{i} g_{i}\left(\bar{x}_{c}(\lambda)\right)<0$ whenever $\sum_{i=1}^{m} \lambda_{i}$ is sufficiently large. In particular, $\exists j \in$ $\{1, \ldots, m\}$ such that $g_{j}\left(\bar{x}_{c}(\lambda)\right)<0$.

Proof: We prove the result by contradiction. Suppose $\sum_{i=1}^{m} \lambda_{i} g_{i}\left(\bar{x}_{c}(\lambda)\right) \geq 0$ whenever $\sum_{i=1}^{m} \lambda_{i}$ is sufficiently large. By optimality of $\bar{x}_{c}(\lambda)$, we must have $L_{c}\left(\bar{x}_{c}(\lambda), \lambda\right) \leq$ $L_{c}(x, \lambda), \forall x \neq \bar{x}_{c}(\lambda)$ and by assumption (A3) $\exists \tilde{x}$ such that $g_{i}(\tilde{x})<0, \forall i \in\{1, \ldots, m\}$. Then $L_{c}\left(\bar{x}_{c}(\lambda), \lambda\right)-L_{c}(\tilde{x}, \lambda)=$ $f\left(\bar{x}_{c}(\lambda)\right)-c^{\top} \bar{x}_{c}(\lambda)+\sum_{i=1}^{m} \lambda_{i} g_{i}\left(\bar{x}_{c}(\lambda)\right)-f(\tilde{x})+c^{\top} \tilde{x}-$ $\sum_{i=1}^{m} \lambda_{i} g_{i}(\tilde{x}) \geq f\left(\bar{x}_{c}(\lambda)\right)-c^{\top} \bar{x}_{c}(\lambda)-f(\tilde{x})+c^{\top} \tilde{x}-$ $\max _{k} g_{k}(\tilde{x}) \sum_{i=1}^{m} \lambda_{i}$.

By assumptions (A1) and (A2) $\tilde{x}$ lies in a compact set, furthermore $f\left(\bar{x}_{c}(\lambda)\right)-c^{\top} \bar{x}_{c}(\lambda)$ and $f(\tilde{x})-c^{\top} \tilde{x}$ are bounded from below, since by (A2) they possess a minimum. We see that $L_{c}\left(\bar{x}_{c}(\lambda), \lambda\right)-L_{c}(\tilde{x}, \lambda)>0$ whenever $\sum_{i=1}^{m} \lambda_{i}$ has been chosen sufficiently large. This contradicts the assumption that $\sum_{i=1}^{m} \lambda_{i} g_{i}\left(\bar{x}_{c}(\lambda)\right) \geq 0$ whenever $\sum_{i=1}^{m} \lambda_{i}$ is large.

Since $\sum_{i=1}^{m} \lambda_{i} g_{i}\left(\bar{x}_{c}(\lambda)\right)<0$ and $\lambda_{i} \geq 0$, we conclude that $\exists j \in\{1, \ldots, m\}$ such that $g_{j}\left(\bar{x}_{c}(\lambda)\right)<0$.
In order to use the previously defined Lyapunov function $V_{1}(x, \lambda)$ in (13), we first show that it is radially unbounded in $x$, i.e. $V_{1}(x, \lambda) \rightarrow \infty$ whenever $\|x\| \rightarrow \infty$. For this purpose, we introduce a set containing all $x$ belonging to the level sets of $V_{1}(x, \lambda)$ and prove that this set is always compact. Let $M_{\alpha}$ denote the set $M_{\alpha}=\left\{x \in \mathbb{R}^{n}\right.$ : $\left.\|\nabla L(x, \lambda)\|=\alpha, \lambda \in \mathbb{R}_{+}^{m}\right\}$.

Lemma 2: Let assumptions (A1) - (A3) be satisfied, then for any fixed $\alpha \geq 0$ the set $M_{\alpha}$ is nonempty and compact. Furthermore, $V_{1}(x, \lambda)$ is radially unbounded in $x$.

Proof: We will first proof that the set

$$
\begin{align*}
N_{c} & =\left\{x \in \mathbb{R}^{n}: x=\bar{x}_{c}(\lambda), \lambda \in \mathbb{R}_{+}^{m}\right\}  \tag{11}\\
& =\left\{x \in \mathbb{R}^{n}: \nabla L(x, \lambda)=c, \lambda \in \mathbb{R}_{+}^{m}\right\}
\end{align*}
$$

is compact. In order to prove that $N_{c}$ is compact, we will make use of the fact that $\bar{x}_{c}(\lambda)$ defined by (10) is unique and exists for all $\lambda \in \mathbb{R}_{+}^{m}$. By assumptions (A1) and (A2), $\bar{x}_{c}=$ $\bar{x}_{c}(\lambda)$ is a well-defined continuous function and therefore
there exists a unique and bounded $\bar{x}_{c}(\lambda)$ for $\sum_{i=1}^{m} \lambda_{i}$ being in a compact set. Hence, it remains to show that $\bar{x}_{c}=\bar{x}_{c}(\lambda)$ is bounded whenever $\sum_{i=1}^{m} \lambda_{i} \rightarrow \infty$. This directly follows from Lemma 1, since for $\sum_{i=1}^{m} \lambda_{i}$ sufficiently large $\exists j \in$ $\{1, \ldots, m\}$ s.t. $g_{j}\left(\bar{x}_{c}(\lambda)\right)<0$ and since the functions $g_{i}$ admit compact level sets due to assumptions (A1) and (A2), $\bar{x}_{c}(\lambda)$ must be bounded.

Observe now that $M_{\alpha}=\bigcup_{\|c\|=\alpha} N_{c}$, since

$$
\begin{align*}
& \tilde{x} \in\left\{x \in \mathbb{R}^{n}:\|\nabla L(x, \lambda)\|=\alpha, \lambda \in \mathbb{R}_{+}^{m}\right\} \\
& \text { if and only if } \\
& \tilde{x} \in \bigcup_{\|c\|=\alpha}\left\{x \in \mathbb{R}^{n}: \nabla L(x, \lambda)=c, \lambda \in \mathbb{R}_{+}^{m}\right\} \tag{12}
\end{align*}
$$

and since $\|c\|=\alpha$ defines a compact set, we conclude that $M_{\alpha}$ is nonempty and compact.

Finally, we prove that $V_{1}(x, \lambda)$ admits compact level sets in $x$. Equivalently to radially unboundedness, i.e. $\|x\| \rightarrow \infty \Rightarrow V_{1}(x, \lambda) \rightarrow \infty$, we can also write $V_{1}(x, \lambda) \leq \alpha \Rightarrow\|x\| \leq \gamma$. The latter follows immediately from the obtained result, since $M_{\alpha}$ denotes the set of $x$ contained in a level set of $V_{1}(x, \lambda)$ and as $M_{\alpha}$ is bounded for all $\lambda \in \mathbb{R}_{+}^{m}$ and any fixed $\alpha \geq 0$, we conclude that $V_{1}(x, \lambda)$ is radially unbounded in $x$.
We are now ready to prove the existence and boundedness of $x(t)$ and $\lambda(t)$ using the Lyapunov function $V_{1}(x, \lambda)$. A crucial issue in the proof is the existence of trajectories, especially of $\lambda(t)$, as $V_{1}(x, \lambda)$ is only radially unbounded in $x$, i.e. a negative derivative of $V_{1}(x, \lambda)$ along the trajectories of (5) does not imply boundedness of $\lambda(t)$ but only of $x(t)$.

Lemma 3: Let assumptions (A1) - (A3) be satisfied. Moreover, suppose a solution of $(x(t), \lambda(t))$ of (5) initialized in $\mathbb{R}^{n} \times \mathbb{R}_{++}^{m}$ is forward complete. Then $\lim _{t \rightarrow \infty}\|\nabla L(x(t), \lambda(t))\| \rightarrow 0$, equivalently $\lim _{t \rightarrow \infty}\|x(t)-\bar{x}(\lambda(t))\|=0$.

Proof: Consider the Lyapunov function $V_{1}(x, \lambda)=$ $\frac{1}{2} \nabla L(x, \lambda)^{\top} \nabla L(x, \lambda)$ that is by Lemma 2 radially unbounded in $x$.
The derivative along the trajectories of (5) yields

$$
\begin{align*}
\dot{V}_{1} & =(\nabla L(x, \lambda))^{\top}\left(\nabla^{2} L(x, \lambda) \dot{x}+\sum_{i=1}^{m} \dot{\lambda}_{i} \nabla g_{i}(x)\right)  \tag{13}\\
& =-\nabla L(x, \lambda)^{\top} \nabla^{2} L(x, \lambda) \nabla L(x, \lambda)
\end{align*}
$$

Since $\nabla^{2} L(x, \lambda)$ is positive definite and $x(t)$ is forward complete by assumption, $x(t)$ converges to $\nabla L(x, \lambda)=0$. We conclude that $x(t)$ approaches $\bar{x}(\lambda(t))$ as $t \rightarrow \infty$ and furthermore $\lim _{t \rightarrow \infty}\|\nabla L(x, \lambda)\| \rightarrow 0$ which implies $\lim _{t \rightarrow \infty} x(t) \in M_{0}=N_{0}$.

Lemma 4: Let assumptions (A1) - (A3) be satisfied. Then any solution $(x(t), \lambda(t))$ of (5) initialized in $\mathbb{R}^{n} \times \mathbb{R}_{++}^{m}$ is forward complete and bounded.

Proof: First, we show that the solution $(x(t), \lambda(t))$ is forward complete. Due to Lemma 3, $x=x(t)$ is bounded
over the maximal interval of existence of $\lambda=\lambda(t)$ since $h(t):=\|\nabla L(x(t), \lambda(t))\|$ is bounded and monotonically decreasing and thus $x(t) \in M_{h(t)}$ is always bounded as long as the solution $\lambda(t)$ exists. Therefore, the solution is forward complete (bounded), if $\lambda=\lambda(t)$ exists for all $t \geq 0$ (is bounded).

The existence of $\lambda=\lambda(t)$ follows by integrating (5b), that yields $\lambda_{i}(t)=\lambda_{i 0} e^{\int_{0}^{t} g_{i}(x(\tau)) d \tau}, i=1, \ldots, m$.

Since $x(\tau)$ stays bounded for $\tau \in[0, t]$, the integral $\int_{0}^{t} g_{i}(x(\tau)) d \tau$ has a finite value for any $t \geq 0$. Hence, $\lambda(t)$ is forward complete.

Finally we show that $\sum_{i=1}^{m} \lambda_{i}(t)$ is bounded. Consider the Lyapunov function $W(\lambda)=\sum_{i=1}^{m} \lambda_{i}$ that is positive definite on $\mathbb{R}_{+}^{m}$ and radially unbounded. The derivative along the trajectories yields $\dot{W}=\sum_{i=1}^{m} \dot{\lambda}_{i}=\sum_{i=1}^{m} \lambda_{i} g_{i}(x)$. Substituting $x=\xi+\bar{x}(\lambda)$, we obtain $\dot{W}=\sum_{i=1}^{m} \lambda_{i} g_{i}(\xi+\bar{x}(\lambda))$. Due to Lemma 1 we have $\sum_{i=1}^{m} \lambda_{i} g_{i}(\bar{x}(\lambda))<0$ whenever $\sum_{i=1}^{m} \lambda_{i}$ is sufficiently large. Moreover, due to Lemma 3 we have $\lim _{t \rightarrow \infty}\|\xi(t)\|=0$, and with forward completeness of $(x(t), \lambda(t))$, we conclude that there exists a time $t_{1}>0$ such that $\forall t>t_{1}$ we have $\dot{W}=\sum_{i=1}^{m} \lambda_{i} g_{i}(\xi+\bar{x}(\lambda))<0$ whenever $\sum_{i=1}^{m} \lambda_{i}$ is sufficiently large. Thus, since $W(\lambda)$ is radially unbounded on $\mathbb{R}_{+}^{m}, \lambda=\lambda(t)$ stays bounded.

As seen before, a crucial issue in this work is to establish forward completeness and boundedness of the solutions. In the proof, the use of assumptions (A2) and (A3) plays an important role. The fact that $\sum_{i=1}^{m} \lambda_{i} g_{i}(\bar{x}(\lambda))$ is strictly smaller than zero assures that there exists a finite time such that $\dot{W}<0$.

We are now ready to state our main result, namely the convergence of $(x(t), \lambda(t))$ to a saddle point $\left(x^{*}, \lambda^{*}\right)$ of $L(x, \lambda)$ in (2) and therefore to the solution of (1).

Theorem 3: Let assumptions (A1) - (A3) be satisfied, suppose $f$ and $g_{i}, i=1, \ldots, m$ are analytic functions and let $\left(x^{*}, \lambda^{*}\right)$ be the unique saddle point of $L$. Then any solution of $(x(t), \lambda(t))$ of (5a) and (5b) initialized in $\mathbb{R}^{n} \times \mathbb{R}_{++}^{m}$ converges to a single connected component $E_{i}$ of the set

$$
\begin{equation*}
E=\left\{(x, \lambda) \in \mathbb{R}^{n} \times \mathbb{R}_{+}^{m}: x=\bar{x}(\lambda), \lambda_{i} g_{i}(x)=0\right\} \tag{14}
\end{equation*}
$$

Furthermore there exist finitely many connected components $E_{i}$, where all of them are unstable except the singleton $E_{0}$ containing the saddle point $\left(x^{*}, \lambda^{*}\right)$ that is asymptotically stable.

Proof: Step 1: We show that all points $\left(\bar{x}^{*}, \bar{\lambda}^{*}\right)$ with $\bar{x}^{*}=\bar{x}\left(\bar{\lambda}^{*}\right)$ and $\bar{\lambda}^{*}=\arg \max _{\lambda \in \mathbb{R}_{+}^{m}} L(\bar{x}(\lambda), \lambda)$ are saddle points. This can be verified using the necessary condition for a function having a maximum at $\bar{\lambda}^{*} \in \mathbb{R}_{+}^{m}$ (see [5] , p. 177, Ex. 1.1):

$$
\begin{align*}
\left.\left(\nabla L(x, \lambda) \frac{\partial \bar{x}(\lambda)}{\partial \lambda}+\frac{\partial L(x, \lambda)}{\partial \lambda}\right)\right|_{\bar{x}\left(\bar{\lambda}^{*}\right)} & =g_{i}\left(\bar{x}\left(\bar{\lambda}^{*}\right)\right) \leq 0 \\
\bar{\lambda}_{i}^{*} g_{i}\left(\bar{x}\left(\bar{\lambda}^{*}\right)\right) & =0, i=1, \ldots, m \tag{15}
\end{align*}
$$

With this and since $x=\bar{x}(\lambda)$ all KKT-Conditions are fulfilled and all $\left(\bar{x}\left(\bar{\lambda}^{*}\right), \bar{\lambda}^{*}\right)$ are saddle points of $L$. Since
$\left(x^{*}, \lambda^{*}\right)$ is unique, the maximum of $L(\bar{x}(\lambda), \lambda)$ is unique and $\left(\bar{x}^{*}, \bar{\lambda}^{*}\right)=\left(x^{*}, \lambda^{*}\right)$.

Step 2: We conclude with assumptions (A1) - (A3) and by Lemma 4 that $x(t)$ and $\lambda(t)$ are bounded. Choose $\mathcal{M}:=\mathbb{R}^{n} \times \mathbb{R}^{m}$ as the manifold with the standard Euclidian metric and define the set $\mathcal{S}=\left\{[x, \lambda]^{\top}: x=\bar{x}(\lambda), \lambda \in \mathbb{R}_{+}^{m}\right\} \cap \Omega$ with $\Omega$ a compact set containing the $\omega$-limit set of $(x(t), \lambda(t))$ initialized at $\mathbb{R}^{n} \times \mathbb{R}_{++}^{m}$. Obviously, we have $\mathcal{S} \subset \mathcal{M}$. By Lemma 3 we conclude that the set $\mathcal{S}$ is attractive for solutions of (5).

We introduce the height (dual) function $V_{2}(\lambda):=$ $-L(\bar{x}(\lambda), \lambda)$ and calculate the derivative $\dot{V}_{2}$ along the trajectories of (5) on the set $\mathcal{S}$. This yields $\dot{V}_{2}=-\left.\left(\nabla L(x, \lambda) \frac{\partial \bar{x}(\lambda)}{\partial \lambda}+\frac{\partial L(x, \lambda)}{\partial \lambda}\right)\right|_{\bar{x}(\lambda)} \dot{\lambda}=$ $-\sum_{i=1}^{m} g_{i}(\bar{x}(\lambda))^{2} \lambda_{i}$.
Define the set $E:=\left\{[x, \lambda]^{\top}: x=\bar{x}(\lambda), \lambda \in \mathbb{R}_{+}^{m}, e(\lambda)=0\right\}$ with $e(\lambda): \lambda \mapsto\left[\lambda_{1} g_{1}^{2}(\bar{x}(\lambda)), \ldots, \lambda_{m} g_{m}^{2}(\bar{x}(\lambda))\right]$. Observe that $E$ is the set where $\dot{V}_{2}=0$, i.e. $\dot{V}_{2}<0$ on $\mathcal{S} \backslash E$ and it contains the set of equilibrium points satisfying the KKT-Conditions (4a), (4b).

Step 3: Define the set $E_{\lambda}:=\left\{\lambda \in \mathbb{R}_{+}^{m}: e(\lambda)=0\right\}$ and notice that every connected component of $E_{\lambda}$ is contained in a level set of $V_{2}(\lambda)$ if and only if every connected component of $E$ is contained in $V_{2}(\lambda)$, since $V_{2}(\lambda)=L(\bar{x}(\lambda), \lambda)$.

In order to apply Theorem 6 in [3] (see Appendix VIB), we will now show that every connected component of $E_{\lambda}$ is contained in a level set of $V_{2}(\lambda)$. In order to prove this, we show that $V_{2}(\lambda)$ is constant in directions of the tangent cone $T_{\lambda}^{C} E_{\lambda}$ at any point $\lambda \in E_{\lambda}$, i.e. $\left.\frac{d}{d \epsilon} V_{2}(\lambda+\epsilon w)\right|_{\epsilon=0}=\frac{\partial V_{2}(\lambda)}{\partial \lambda} w=0, \forall w \in T_{\lambda}^{C} E_{\lambda}$ with $\frac{\partial V_{2}(\lambda)}{\partial \lambda}=\left.\frac{\partial L(x, \lambda)}{\partial \lambda}\right|_{\bar{x}(\lambda)}=\left[g_{1}(\bar{x}(\lambda)), \ldots, g_{m}(\bar{x}(\lambda)]\right.$, because $\left.\nabla L(x, \lambda)\right|_{\bar{x}(\lambda)}=0$. Since (see e.g. [12, p. 44]) $T_{\lambda}^{C} E_{\lambda} \subseteq$ $R_{\lambda} E_{\lambda}=\left\{v \in \mathbb{R}^{m}:\left.\frac{\partial e}{\partial \lambda}\right|_{\lambda \in E_{\lambda}} v=0\right\}$, it suffices to prove that $\frac{\partial V_{2}(\lambda)}{\partial \lambda} v=0, \forall v \in R_{\lambda} E_{\lambda}$. This implies that we stay on a level set of $V_{2}(\lambda)$ for a small perturbation in the direction of $v \in R_{\lambda} E_{\lambda}$, and therefore also for any $w \in T_{\lambda}^{C} E_{\lambda}$. The set $E_{\lambda}$ contains all points with $\lambda_{i} g_{i}^{2}(\bar{x}(\lambda))=0$. We distinguish between two cases:

1) $\forall i \in\{1, \ldots, m\}: g_{i}(\bar{x}(\lambda))=0$.

It follows that $\frac{\partial V_{2}(\lambda)}{\partial \lambda}=0$. Since the dual function is concave $\left(V_{2}(\lambda)\right.$ is convex, see [5]) each point where the gradient vanishes is a global maximizer of the dual function, and consequently the connected components where the gradient vanishes, must all lie on the same level set ${ }^{1}$.
2) $\exists i \in\{1, \ldots, m\}: g_{i}(\bar{x}(\lambda)) \neq 0$.

Without loss of generality, we relabel all $g_{i}(\bar{x}(\lambda))$ such that:

$$
\begin{aligned}
& \text { - } g_{i}(\bar{x}(\lambda)) \neq 0, i=1, \ldots, k \\
& \text { - } g_{i}(\bar{x}(\lambda))=0, i=k+1, \ldots, m
\end{aligned}
$$

with $1 \leq k \leq m$. Then, the Jacobian of $e(\lambda)$ at $\lambda \in E$

[^1]yields:
\[

\left.\frac{\partial e}{\partial \lambda}\right|_{\lambda \in E}=\left[$$
\begin{array}{cc}
\operatorname{diag}\left[g_{1}^{2}(\bar{x}(\lambda)), \ldots, g_{k}^{2}(\bar{x}(\lambda))\right] & 0  \tag{16}\\
0 & 0
\end{array}
$$\right]
\]

Therefore, all vectors in $R_{\lambda} E_{\lambda}$ must be of the form:

$$
\begin{equation*}
v=\left[0, \ldots, 0, v_{k+1}, \ldots, v_{m}\right]^{\top} \tag{17}
\end{equation*}
$$

and the directional derivative of $V_{2}(\lambda)$ along the elements in $R_{\lambda} E_{\lambda}$ yields $\frac{\partial V_{2}(\lambda)}{\partial \lambda} v=\sum_{i=k+1}^{m} g_{i}(\bar{x}(\lambda)) v_{i}=$ $0, \forall v \in R_{\lambda} E_{\lambda}$.
Therefore, all vectors $v \in R_{\lambda} E_{\lambda}$ are perpendicular to the gradient on the level set of $V_{2}(\lambda)$ and since $T_{\lambda}^{C} E \subseteq R_{\lambda} E_{\lambda}$ all points of a connected component of $E$ lie in a level set of $V_{2}(\lambda)$. Thus, $E_{\lambda}$ is contained in a level set of $V_{2}(\lambda)$.
The set $E$ can be decomposed such that $E=\bigcup_{i \in I} E_{i}$ with $\left\{E_{i}\right\}_{i \in I}$ denoting the connected components in $E$ and let $\nu^{i} \in \mathcal{V}=V_{2}(E)$ denote the image of of $E$ under $V_{2}(\lambda)$ (we write $V_{2}(E)$ instead of $-L(E)$ ), i.e. $\nu^{i}=V_{2}\left(E_{i}\right)$. Since all $\left\{E_{i}\right\}_{i \in I}$ are contained in the level sets of $V_{2}(\lambda)$, we conclude that different $\nu^{i}$ 's must posses different preimages with $z^{i}=\left[x^{i}, \lambda^{i}\right]^{\top} \in E_{i}$. Remember that $x^{i}$ is uniquely defined for a fixed $\lambda^{i}$, since $z^{i} \in E_{i}$ with $x=\bar{x}(\lambda)$, therefore $E=\left\{[x, \lambda]^{\top}\right.$ : $\left.\nabla f(x)+\sum_{i=1}^{m} \lambda_{i} \nabla g_{i}(x)=0, \lambda_{i} g_{i}(x)=0, i=1, \ldots, m\right\}$, and thus $z^{i}$ is a zero of an analytic function. We will now exclude the existence of accumulation points in $\mathcal{V}$. Suppose, there exists an accumulation point $\nu^{\infty}$, i.e. there exists a sequence $\left\{\nu^{k}\right\}$, with distinct $\nu^{k}$ such that (see [19]) $\lim _{k \rightarrow \infty} \nu^{k}=\nu^{\infty}$. Let $z^{k}$ denote some preimage of $\nu^{k}$, i.e. $z^{k} \in E_{k}$. The set $E$ is compact, so there exists a converging subsequence $\left\{\tilde{z}^{j}\right\} \subseteq\left\{z^{k}\right\}$ (Thm. of BolzanoWeierstrass, [19]) with distinct $\tilde{z}^{j}$ such that $\lim _{j \rightarrow \infty} \tilde{z}^{j}=$ $\tilde{z}^{\infty}, \tilde{z}^{\infty} \in V_{2}^{-1}\left(\nu^{\infty}\right)$. Since $\tilde{z}^{\infty}$ is the accumulation point of $\left\{\tilde{z}^{j}\right\}$, there exists for every $\epsilon>0$ a $p \in \mathbb{N}$ such that $\tilde{z}^{j} \in B_{\epsilon}\left(\tilde{z}^{\infty}\right), \forall j>p$, i.e. for any $\epsilon>0$ there exist infinitely many distinct $\tilde{z}^{j} \in B_{\epsilon}\left(\tilde{z}^{\infty}\right)$ and therefore infinitely many distinct connected components $E_{j}\left(\tilde{z}^{j} \in E_{j}\right)$ intersect with an $\epsilon$-ball $B_{\epsilon}\left(\tilde{z}^{\infty}\right)$. This leads to a contradiction since $f$ and $g_{i}, i=1, \ldots, m$ are analytic and not identically zero (there exists an $x$ and $d$ such hat $h(\alpha)=f(x+\alpha d) \not \equiv 0)$ it follows by Lojasiewicz's Theorem (see Thm. 6.3.3 in [18]) that the zero set of a real valued analytic functions is locally a union of a finite number of connected components. This contradicts the assumption that an accumulation point $\tilde{z}^{\infty}$ exists. We conclude that the connected components in $E$ are isolated and are embedded in the level sets of $V_{2}(\lambda)$.

We are now ready to apply Theorem 6 in [3] (see Appendix VI-B), concluding that the trajectories of (5) approach a connected component $E_{i}$ of $E$ with $t \rightarrow \infty$.

Observe furthermore that $E$ contains all equilibrium points of (5) which satisfy (4a) and (4b) but not necessarily (4c).

Step 4: We will now show by using Chetaev's Instability Theorem that all $E_{i}$ are unstable except the singleton $E_{0}$ containing the saddle point $\left(x^{*}, \lambda^{*}\right)$. Take again $V_{2}(\lambda)=-L(\bar{x}(\lambda), \lambda)$ and define $W_{i}(\lambda)=-V_{2}(\lambda)+c_{i}$ with $\left(x^{i}, \lambda^{i}\right) \in E_{i}, i \geq 1$ and the constant $c_{i}=L\left(x^{i}, \lambda^{i}\right)$. Obviously $W_{i}\left(\lambda^{i}\right)=0$. Define $\bar{B}_{\epsilon}:=B_{\epsilon}\left(x^{i}, \lambda^{i}\right) \cap N_{0} \times \mathbb{R}_{+}^{m}$,
$\left(x^{i}, \lambda^{i}\right) \in E_{i}, i \geq 1$ with $N_{0}$ as in (11). Since $W_{i}\left(\lambda^{*}\right)>$ $W_{i}\left(\lambda^{i}\right), i \geq 1$ and $W_{i}, i \geq 0$ are concave, it follows that for every $\epsilon$ there exists a subset $D$ of $\bar{B}_{\epsilon}$ such that $W_{i}(\lambda)>0$ for all $\lambda \in D \backslash \lambda^{i}$. Now since $N_{0}$ is invariant and $\dot{W}_{i}=\sum_{i=1}^{n} g_{i}(\bar{x}(\lambda))^{2} \lambda_{i}>0$, for all $(\bar{x}(\lambda), \lambda) \in \bar{B}_{\epsilon}$ (and therefore also on $D$ ) it follows by Chetaev's Instability Theorem (see [16]) that the sets $E_{i}, i \geq 1$ are unstable.

In the following we use Theorem 11 in [15] in order to show that $E_{0}$, i.e. the saddle point, is asymptotically stable. Note that $V_{1}$ in (13) is a positive semi-definite Lyapunov function on $\mathbb{R}^{n} \times \mathbb{R}_{+}^{m}$ that vanishes on the set $N_{0}$. In order to show that the saddle-point $\left(x^{*}, \lambda^{*}\right)$ is asymptotically stable it is sufficient to show that the restriction of (5) is asymptotically stable on the set $N_{0}{ }^{2}$. Using the Lyapunov function $-W_{0}$ that is positive definite on $N_{0}$ and vanishes only at the saddle-point, the Lie-derivative yields $\dot{W}_{0}=$ $-\sum_{i=1}^{n} g_{i}(\bar{x}(\lambda))^{2} \lambda_{i}$ and is strictly negative definite for all $\lambda \neq \lambda^{*}$. This proves that the saddle point $\left(x^{*}, \lambda^{*}\right) \in E_{0}$ is asymptotically stable on $N_{0}$. Using the results of Theorem 11 in [15] we can conclude from the asymptotic stability of $\left(x^{*}, \lambda^{*}\right)$ on $N_{0}$ that $\left(x^{*}, \lambda^{*}\right)$ is asymptotically stable.
In the following, we discuss the main result and the assumption under which it has been established as well as the applicability of (5) to linear programming.

First note that the existence of a unique saddle point of $L$ is assumed to be able to use Theorem 11 in [15], since it is formulated such that it is valid for a point. It can be expected however that this is also valid in the case that $E_{0}$ is not a singleton.

Second it is important to point out that the system converges to $E_{0}$ for practically every initial condition. The convergence to an unstable connected component $E_{i}$ is highly unlikely, but possible in some pathological cases.

We restrict the functions $f$ and $g_{i}, i=1, \ldots, m$ to be analytic because a finite number of accumulation points in $V_{2}(E)$ is required for Theorem 6 in [3]. Since analytic functions can have only isolated connected components, we assure the non-existence of accumulation points. On the other hand, infinitely many accumulation points can only occur in highly degenerated cases. As pointed out in [11], most nonlinear differentiable programs are well-behaved. Furthermore, a finite number of accumulation points still leads to a guaranteed convergence but there is no general criteria known to the authors that gives a reasonable characterization of such functions.

Finally, the assumptions (A1) - (A3) are conservative but needed in our current proof, mainly in order to establish boundedness of solutions. From a practical point of view, one might expect that (5) will converge in more general situations when the assumptions are not satisfied.

The application of (5) can be extended to linear programs. The proof in this special case can be established using a single Lyapunov-function and under less strict assumptions as those needed in Theorem 3.

[^2]

Fig. 2: Feasible Sets and Trajectory $\left(x_{1}(t), x_{2}(t)\right)$

In the following we will present the structure of the dynamical system for linear programs.

First recall a linear program in inequality form:

$$
\begin{align*}
& \inf _{x} c^{\top} x  \tag{18}\\
& \text { s.t. } a_{i}^{\top} x-b_{i} \leq 0, i=1, \ldots, m
\end{align*}
$$

We propose the following vector field for this type of problems:

$$
\begin{align*}
\dot{x} & =-c-\sum_{i=1}^{m} \lambda_{i} a_{i}\left(a_{i}^{\top} x-b_{i}+1\right)  \tag{19a}\\
\dot{\lambda}_{i} & =\lambda_{i}\left(a_{i}^{\top} x-b_{i}\right), i=1, \ldots, m \tag{19b}
\end{align*}
$$

This system is a special realization of (5) for linear programs where $\nabla^{2} L(x, \lambda)^{-1}$ was substituted by the identity matrix. Under the assumption that the saddle point $\left(x^{*}, \lambda^{*}\right)$ of $L(x, \lambda)=c^{\top} x+\sum_{i=1}^{n} \lambda_{i}\left(a_{i}^{\top} x-b_{i}\right)$ is unique, it can be shown that the trajectories of (19) initialized in $\mathbb{R}^{n} \times \mathbb{R}_{++}^{m}$ converge to $\left(x^{*}, \lambda^{*}\right)$ and that it is exponentially stable. Thus, the system converges to the solution $x^{*}$ of (18). The proof can be found in [8].

## IV. Example

In the following we show two examples. In the first example we apply the algorithm (5) to the optimization problem

$$
\begin{align*}
& \min _{x}\left(x_{1}-1\right)^{2}+\left(x_{2}-2\right)^{2} \\
& \text { s.t. }\left(x_{1}-1\right)^{2}+x_{2}^{2}-2 \leq 0  \tag{20}\\
& \quad\left(x_{1}+1\right)^{2}+x_{2}^{2}-2 \leq 0
\end{align*}
$$

The trajectory of the $x$-coordinate as well as the feasible set and the level sets of the objective function are shown in Fig. 2. One can also see the trajectory of $\bar{x}(\lambda(t))$ and its tracking by $x(t)$.

We emphasize, that there is no need to choose an initial condition in the feasible set, but only positive initial conditions $\lambda_{i 0}$.

## V. Summary and Future Work

We proposed the vector field (5) whose trajectories converge globally to a saddle point of the Lagrangian function (2) belonging to the constrained optimization problem (1).

Under suitable assumptions such as strict convexity and strict feasibility it is well known that the saddle point of the Lagrangian coincides with the solution of a convex optimization problem. Thus, (5) provides a novel way to solve convex optimization problems of type (1).

One nice feature of system (5) is that the well-known regularity and optimality conditions from optimization theory do not only guarantee well-posedness of the problem but also dictate local and global qualitative behavior of the vector field (5). The main focus of our future research lies in the relaxation of the conditions on the functions $f$ and $g_{i}, i=1, \ldots, m$, such that we are able to prove global convergence for other classes of optimization problems.

Another interesting research point is to study the underlying geometry of (5). Saddle point problems possess interesting geometric properties (see [6] and [11]). As shown in [6], saddle point algorithms like the AHU-flow have a gradient-like structure using a suitable (indefinite) metric. It is interesting in that respect that the $\lambda$-dynamics of (5) can be interpreted as a gradient flow for the dual function with the definite (but singular) metric $\operatorname{diag}\left[\lambda_{1}^{-1}, \ldots, \lambda_{m}^{-1}\right]$.

## VI. APPENDIX

## A. Derivative of $\bar{x}_{c}(\lambda)$ with respect to $\lambda$

Lemma 5: Suppose (A1) and (A2) are satisfied and let $\bar{x}_{c}(\lambda)$ denote the unique minimizer of $L_{c}(x, \lambda)$ for a given $\lambda \in \mathbb{R}_{+}^{m}$. Then $\frac{\partial \bar{x}_{c}(\lambda)}{\partial \lambda}=$ $-\left(\nabla^{2} L\left(\bar{x}_{c}(\lambda), \lambda\right)\right)^{-1}\left[\nabla g_{1}\left(\bar{x}_{c}(\lambda)\right), \ldots, \nabla g_{m}\left(\bar{x}_{c}(\lambda)\right)\right]$. By abuse of notation, we write $\nabla^{2} L\left(\bar{x}_{c}(\lambda), \lambda\right)$ instead of $\left.\frac{\partial \nabla L(x, \lambda)}{\partial x}\right|_{\bar{x}_{c}(\lambda)}$.
Proof: The minimizer $\bar{x}_{c}(\lambda)$ is implicitly defined as the unique solution of the equation (10). Due to (A1) and (A2), the minimizer exists and is unique. Moreover, due to the Implicit Function Theorem, the derivative of $\bar{x}_{c}(\lambda)$ with respect to $\lambda$ exists and can be calculated by differentiating both sides of (10) with respect to $\lambda$

$$
\begin{equation*}
\left.\frac{\partial \nabla L(x, \lambda)}{\partial x}\right|_{\bar{x}_{c}(\lambda)} \frac{\partial \bar{x}_{c}(\lambda)}{\partial \lambda}+\left.\frac{\partial \nabla L(x, \lambda)}{\partial \lambda}\right|_{\bar{x}_{c}(\lambda)}=0 \tag{21}
\end{equation*}
$$

where $\left.\frac{\partial \nabla L(x, \lambda)}{\partial x}\right|_{\bar{x}_{c}(\lambda)}$ coincides with $\left.\nabla^{2} L(x, \lambda)\right|_{\bar{x}_{c}(\lambda)}$. This yields $\frac{\partial \bar{x}_{c}(\lambda)}{\partial \lambda}=-\left.\left(\nabla^{2} L\left(\bar{x}_{c}(\lambda), \lambda\right)\right)^{-1} \frac{\partial \nabla L(x, \lambda)}{\partial \lambda}\right|_{\bar{x}_{c}(\lambda)}=$ $-\left(\nabla^{2} L\left(\bar{x}_{c}(\lambda), \lambda\right)\right)^{-1}\left[\nabla g_{1}\left(\bar{x}_{c}(\lambda)\right), \ldots, \nabla g_{m}\left(\bar{x}_{c}(\lambda)\right)\right]$.

## B. Theorem 6 in [3]

The following assumptions are needed for the Theorem: 1. Consider a manifold $\mathcal{M} \subseteq \mathbb{R}^{N}$ of class $C^{2}$ with standard metric on which a locally Lipschitz continuous vector field $\dot{x}=f(x)$ is given; 2 . The solution $x(t)$ is bounded; 3. The $\omega$-limit set $\Omega(x(0))$, which is a compact and connected set, is contained in a closed embedded submanifold $\mathcal{S} \subset \mathcal{M}$. Equivalently, assume that $\mathcal{S}$ is attracting for $x(t)$ starting at $x(0) ; 4$. Call $O$ an open tubular neighborhood of $\mathcal{S}$ in $\mathcal{M}$ and assume that there exists a real-valued $C^{1}$ height function $W: O \rightarrow \mathbb{R}$ and such that $\dot{W} \geq$ on S (or $\dot{W} \leq 0$ on $\mathcal{S}$ ), where $\dot{W}$ is the derivative of $W(x)$ along the trajectory (Lie derivative). Moreover, let $E:=\{x \in \mathcal{S}: \dot{W}=0\}$ so that $\dot{W}>0$ on $\mathcal{S} \backslash E$ (or $\dot{W}=0$ on $\mathcal{S} \backslash E$ ).

Definition: Let $\left\{E_{i}\right\}_{i \in I}$ be the connected components of $E$. Given a height function $W(x)$ as in the assumptions above, say that the components $\left\{E_{i}\right\}_{i \in I}$ are contained in $W(x)$ if each $E_{i}$ lies in a level set of $W(x)$, and the subset $\left\{W\left(E_{i}\right)\right\}_{i \in I} \subset \mathbb{R}$ has at most a finite number of accumulation points in $\mathbb{R}$.

Theorem 4: Assume the assumptions above hold. If the components $\left\{E_{i}\right\}_{i \in I}$ are contained in $W(x)$ according to the foregoing Definition, then $\Omega(x(0)) \subset E_{i}$ for a unique $i \in I$.

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[^1]:    ${ }^{1}$ The convexity condition is needed, since there exists (non-convex and non-concave) functions not constant on a connected set of points where its gradient vanishes (see [22]).

[^2]:    ${ }^{2}$ In terms of Theorem 11 in [15]: $V_{1}=V$ and $N_{0}=L$ with $L$ being the largest invariant set contained in $\{x: V(x)=0\}$.

