# Optimization Models and Applications 

Edited by<br>Siamak Pedrammehr and Mohammad Reza Chalak Qazani

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Editors

## Siamak Pedrammehr

Mohammad Reza Chalak Qazani

Editors

| Siamak Pedrammehr | Mohammad Reza Chalak Qazani |
| :--- | :--- |
| Faculty of Design | Faculty of Computing and |
| Tabriz Islamic Art University | Information Technology (FoCIT) |
| Tabriz | Sohar University |
| Iran | Sohar |
|  | Oman |

Editorial Office
MDPI
St. Alban-Anlage 66
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## About the Editors

## Siamak Pedrammehr

Siamak Pedrammehr earned his Bachelor's and Master's degrees in Mechanical Engineering from the University of Tabriz in Iran. He further advanced his academic credentials by obtaining a Ph.D. in Engineering from Deakin University in Australia, where he subsequently received the prestigious Alfred Deakin Postdoctoral Research Fellowship in 2019. Following his Ph.D., he was appointed as a postdoctoral research fellow by Iran's National Elites Foundation at the University of Tabriz, a role he fulfilled for two years.

Currently, Dr. Pedrammehr holds a position as an Assistant Professor at the Faculty of Design at Tabriz Islamic Art University. He maintains an active research partnership with Deakin University. His research interests are diverse and include modeling and simulation, manufacturing processes, applied AI and socio-industrial studies.

## Mohammad Reza Chalak Qazani

Mohammad Reza Chalak Qazani commenced his academic journey in manufacturing and production engineering, obtaining his Bachelor's degree from the University of Tabriz, Iran. He continued his studies in robotic and mechanical engineering, earning a master's degree from Tarbiat Modares University, Iran. Furthering his scholarly pursuit, he achieved his Ph.D. from the Institute for Intelligent Systems Research and Innovation (IISRI) at Deakin University, Australia. Upon completion of his doctorate, Dr. Qazani was granted the esteemed Alfred Deakin Postdoctoral Research Fellowship at the IISRI, dedicating two years to the advancement of model predictive control, motion cueing algorithms and soft computing controllers.

Currently, Dr. Qazani serves as an Assistant Professor within the Faculty of Computing and Information Technology (FoCIT) at Sohar University in Oman. His instructional expertise encompasses a spectrum of subjects, including data structures and algorithms, enterprise resource planning modeling and implementation, modeling and visualization, computer architecture, introductory artificial intelligence and advanced machine learning.

## Preface

With its roots in diverse fields, optimization has evolved into an indispensable tool shaping the landscapes of engineering, economics, the environment, health, systems, businesses and more. It has become a guiding light, leading us toward optimal solutions in the intricate web of scientific challenges. In the celestial expanse of optimization, this book explores the convergence of theory and practice, shedding light on the transformative applications across various domains. Furthermore, it serves as a platform bridging theoretical rigor with practical significance, addressing scientific challenges across disciplines such as functional analysis, critical point theory, bifurcation theory, set-valued analysis, calculus of variations, etc. The convergence of theory and reality within this realm aims to advance optimization solutions and provide insights into the fundamental scientific challenges of our time. We extend our heartfelt gratitude to the authors for their outstanding contributions and to the diligent reviewers whose valuable comments have enhanced the quality of the articles. Special thanks to the Editorial Board and the Editorial Office of Axioms for their invaluable support and guidance throughout the publication process. We trust that our readers will find the articles in this journal to be informative, insightful and rich sources of new and valuable information on optimization models and their practical applications.

We wish you a happy reading!

Editorial

# Special Issue "Optimisation Models and Applications" 

Siamak Pedrammehr 1,2,*(D) and Mohammad Reza Chalak Qazani 1,3(D)<br>1 IISRI, Deakin University, Waurn Ponds, VIC 3216, Australia; m.r.chalakqazani@gmail.com<br>2 Faculty of Design, Tabriz Islamic Art University, Tabriz 5164736931, Iran<br>3 FoCIT, Sohar University, Sohar 311, Oman<br>* Correspondence: s.pedrammehr@tabriziau.ac.ir

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[^0]Optimisation models have transcended their origins to become indispensable tools across many fields, including engineering, economics, the environment, health, systems of systems, businesses, and beyond. These models serve as guiding lights, illuminating the path toward optimal solutions. Within the optimisation realm, the following four distinct categories emerge, akin to constellations in the scientific cosmos: physics-based optimisation algorithms, swarm-based optimisation algorithms, game-based optimisation algorithms, and evolutionary-based optimisation algorithms. Over the past half-century, the optimisation domain has been a crucible where theory meets practice, generating solutions that have left indelible marks on diverse applications. In this Special Issue, our study takes us to the heart of the intricate interplay between optimisation models and their practical applications. In this realm, theory and reality merge to address fundamental scientific challenges. We earnestly invite researchers, akin to celestial navigators, to contribute their original, high-quality research papers. Let us explore the celestial expanse of optimisation and its transformative applications.

This Special Issue explores the intersection of optimisation models and real-world applications, encompassing diverse disciplines such as functional analysis, the critical point theory, bifurcation theory, set-valued analysis, calculus of variations, partial differential equations, variational and topological methods, fixed-point theory, game theory, convex analysis, matrix theory, control theory, and data mining. As a platform for researchers, it aims to bridge theoretical rigour with practical significance in addressing scientific challenges and advancing optimisation solutions across domains.

Rasouli et al. (Contribution 1) devised an adaptive sliding mode control for fractionalorder chaotic systems, addressing an unknown time delay, uncertainty, and disturbances. Utilising a nonlinear fractional order PID sliding surface, stability is ensured using the Lyapunov theory and validated through MATLAB simulations. This approach extends to secure image communication, showcasing positive outcomes in correlation with NPCR, PSNR, and information entropy, even under uncertain parameters.

Ma and Xu (Contribution 2) introduced a method tackling robust multi-criteria traffic network equilibrium problems with path-capacity constraints. This study formulates an equivalent min-max optimisation problem, deploying a direct search algorithm and a smoothing optimisation approach based on a ReLU activation function variant for efficient solutions. Numerical examples highlight algorithm efficiency, particularly in small-scale traffic networks.

Giuffrè (Contribution 3) explored a nonconstant gradient-constrained problem, establishing a connection with a double obstacle problem. This study fulfils a constraint qualification condition, enabling the application of solid duality theory. It provides evidence for the existence of Lagrange multipliers, suggesting potential avenues for future research within this framework.

Swarup et al. (Contribution 4) delved into K-bonvexity/K-pseudo convexity concepts, establishing duality theorems for a K-Wolfe multi-objective second-order symmetric duality model. This work broadens mathematical perspectives, potentially contributing to
applications in higher-order symmetric fractional programming problems and variational control problems over cones.

Gismelseed et al. (Contribution 5) employed a mathematical model to investigate factors influencing the sit-to-stand (STS) motion, formulating an optimisation problem to minimise joint torques. The model successfully predicted key STS motion characteristics, considering the motion speed, reduced joint strength, and seat height. These results align with experimental findings, providing insights into STS biomechanics for clinical investigations and daily activities related to human mobility.

Acosta-Portilla et al. (Contribution 6) introduced a methodology for solving the Monge-Kantorovich mass transfer problem, utilising Haar multiresolution analysis and wavelet transform. This approach significantly reduces computational operations, achieving efficient solutions at different resolution levels. Comparative results demonstrate consistent improvement over previous resolution levels, with exact solutions in some cases. This work offers valuable insights and suggests research opportunities for different discretisation techniques and cost function characteristics.

Rajebi et al. (Contribution 7) presented a robust license plate recognition system that is effective in adverse environmental conditions. Unlike prior approaches, this method evaluates and recognises license plates from various sources, handling environmental challenges by removing image artifacts before recognition. Utilising Hopfield's neural network for recognition reduces the execution time. It enhances accuracy compared to traditional methods, contributing significantly to automated surveillance systems.

Pourmostaghimi et al. (Contribution 8) introduced a helical gear reverse engineering methodology employing swarm-based optimisation, specifically the Grey Wolf Optimisation and Particle Swarm Optimisation. Results showcase superior performance in accuracy, convergence speed, and stability. This study suggests diverse gear applications and explores the impact of algorithm tuners on convergence speed.

Rasouli et al. (Contribution 9) proposed a robust synchronisation method for chaotic systems amid uncertainty. Combining the sliding mode control with adaptive rules ensures the convergence of unknown parameters and time delays to zero. Simulations show the robust synchronisation of uncertain, jerk-crusty systems, and the control strategy is applied successfully to secure communication.

Yaghoubi et al. (Contribution 10) present a novel approach for robust stability in polynomial fractional differential (PFD) systems using Caputo derivatives. Employing the sum of squares (SOSs) method, this research addresses stability challenges beyond linear matrix inequalities. Demonstrating robust Mittag-Leffler stability conditions, this study introduces a robust controller for PFD systems with unknown parameters and a polynomial state feedback controller for PFD-controlled systems, validated through simulations, offering innovative solutions for stability and control challenges.

Guo (Contribution 11) investigated a Mond-Weir-type robust duality for uncertain semi-infinite multi-objective fractional optimisation problems. It establishes a robust dual problem, incorporating a new subdifferential constraint qualification and a generalised convex-inclusion assumption. This study unveils robust \#-quasi-weak, strong duality properties, extending previous results. These findings contribute to understanding uncertain fractional optimisation and suggest the future exploration of mixed-type robust approximate dual problems. Specific research funding sources for this work are not provided in the article.

We extend our heartfelt gratitude to the authors for their outstanding contributions to this journal in the field of 'Optimisation Models and Applications'. Our sincere appreciation goes to the diligent reviewers whose valuable comments and feedback have significantly enhanced the quality of the articles. We also acknowledge the invaluable support and guidance provided by the editorial board and the editorial office of Axioms throughout the publication process. We trust that our readers will find the articles in this journal to be informative, insightful, and rich sources of new and valuable information on optimization models and their practical applications.

Author Contributions: All authors contributed equally to writing, editing, and reviewing the manuscript. All authors have read and agreed to the published version of the manuscript.

Conflicts of Interest: The authors declare no conflicts of interest.

## List of Contributions

1. Rasouli, M.; Zare, A.; Yaghoubi, H.; Alizadehsani, R. Designing a Secure Mechanism for Image Transferring System Based on Uncertain Fractional Order Chaotic Systems and NLFPID Sliding Mode Controller. Axioms 2023, 12, 828.
2. Ma, X.-X.; Xu, Y.-D. Robust Multi-Criteria Traffic Network Equilibrium Problems with Path Capacity Constraints. Axioms 2023, 12, 662.
3. Giuffrè, S. A Nonconstant Gradient Constrained Problem for Nonlinear Monotone Operators. Axioms 2023, 12, 605.
4. Swarup, C.; Kumar, R.; Dubey, R.; Fathima, D. New Class of KG-Type Symmetric Second Order Vector Optimization Problem. Axioms 2023, 12, 571.
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6. Acosta-Portilla, J.R.; González-Flores, C.; López-Martínez, R.R.; Sánchez-Nungaray, A. Efficient Method to Solve the Monge-Kantarovich Problem Using Wavelet Analysis. Axioms 2023, 12, 555.
7. Rajebi, S.; Pedrammehr, S.; Mohajerpoor, R. A License Plate Recognition System with Robustness against Adverse Environmental Conditions Using Hopfield's Neural Network. Axioms 2023, 12, 424.
8. Pourmostaghimi, V.; Heidari, F.; Khalilpourazary, S.; Qazani, M.R.C. Application of Evolutionary Optimisation Techniques in Reverse Engineering of Helical Gears: An Applied Study. Axioms 2023, 12, 252.
9. Rasouli, M.; Zare, A.; Hallaji, M.; Alizadehsani, R. The synchronisation of a class of time-delayed chaotic systems using sliding mode control based on a fractional-order nonlinear PID sliding surface and its application in secure communication. Axioms 2022, 11, 738.
10. Yaghoubi, H.; Zare, A.; Alizadehsani, R. Analysis and Design of Robust Controller for Polynomial Fractional Differential Systems Using Sum of Squares. Axioms 2022, 11, 623.
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# The Synchronization of a Class of Time-Delayed Chaotic Systems Using Sliding Mode Control Based on a Fractional-Order Nonlinear PID Sliding Surface and Its Application in Secure Communication 

Mohammad Rasouli ${ }^{1}$, Assef Zare ${ }^{1, *(\mathbb{D}}$, Majid Hallaji ${ }^{2}$ and Roohallah Alizadehsani ${ }^{3}$ (D)<br>1 Faculty of Electrical Engineering, Gonabad Branch, Islamic Azad University, Gonabad 6518115743, Iran<br>2 Faculty of Electrical Engineering, Neyshabure Branch, Islamic Azad University, Neyshabure 6518115743, Iran<br>3 Institute for Intelligent Systems Research and Innovation (IISRI), Deakin University, Geelong, VIC 3216, Australia<br>* Correspondence: assefzare@gmail.com

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#### Abstract

A novel approach for the synchronization of a class of chaotic systems with uncertainty, unknown time delays, and external disturbances is presented. The control method given here is expressed by combining sliding mode control approaches with adaptive rules. A sliding surface of fractional order has been developed to construct the control strategy of the abovementioned sliding mode by employing the structure of nonlinear fractional PID (NLPID) controllers. The suggested control mechanism using Lyapunov's theorem developed robust adaptive rules in such a way that the estimation error of the system's unknown parameters and time delays tends to be zero. Furthermore, the proposed robust control approach's stability has been demonstrated using Lyapunov stability criteria and Lipschitz conditions. Then, in order to assess the performance of the proposed mechanism, the presented control approach was used to simulate the synchronization of two chaotic jerk systems with uncertainty, unknown time delays, and external distortion. The results of the simulation confirm the robust and desirable synchronization performance. Finally, a secure communications mechanism based on the proposed technique is shown as a practical implementation of the introduced control strategy, in which the message signal is disguised in the transmitter with high security and well recovered in the receiver with high quality, according to the mean squared error (MES) criteria.


Keywords: chaotic synchronization; sliding mode control; adaptive control; uncertainty; unknown time-delay; secure communication

MSC: 93D09; 93B51

## 1. Introduction

Chaos is a phenomenon that occurs in nonlinear dynamic systems. The dynamic behavior of these systems is fully dependent on the initial conditions; therefore, even the tiniest change in these parameters generates big changes in their behavior. Many domains of study, including economics [1,2], chemistry [3,4], biology [5,6], and engineering [7,8], have discovered chaotic systems. Many scientists from numerous domains have been interested in the synchronization of chaotic systems in recent years [5]. Pecora and Carroll proposed chaos synchronization in 1990 [6].

Fractional order (FO) controllers were developed by combining fractional calculation with existing controllers. The increased number of configurable parameters in this type of controller creates more flexibility in the control process, which improves the performance of the controlled system. In this regard, the fractional PI controller [7,8], the fractional PD controller [9], the fractional lag-lead controller [10], the fractional CRONE controller [11,12],
the adaptive FO PID controller [13], the reference adaptive control [14], the fractional model [15,16], and the fractional sliding mode control [17-19] have all been mentioned.

Adaptive sliding mode controllers have been shown to exhibit robust performance against system parameter changes and system uncertainties [20,21], with their greatest advantage being stability against system disturbances.

Adaptive control [22-24] is used to achieve control objectives when there is uncertainty in the system. Because of the benefits of stability against system parameter uncertainties, adaptive control has a faster convergence speed.

Sliding mode control, on the other hand, is a popular and successful strategy that is simple to implement $[25,26]$. The sliding mode controller is a variable structure controller that works well for nonlinear systems with model uncertainty. The primary aspect of sliding mode control is that it directs the system's states from the initial states to a suitable sliding surface that is provided and then maintains the states at the mentioned sliding surface for all subsequent iterations. Several articles have been provided in this area for the synchronization of chaotic systems utilizing sliding mode control. Qamati et al. [27] successfully synchronized identical Genesio-Tesi chaotic systems using adaptive sliding mode control. The integrated sliding mode control method was used to examine the control and synchronization of extremely chaotic Zhou systems, as discussed in [28].

The authors of [29] proposed a hybrid fractional-order sliding mode control method for finite-time synchronization of a chaotic class using the direct Lyapunov method. Synchronization of fractional order chaotic systems is introduced in [30]. In this design, a sliding surface based on non-linear fractional order PID is presented. In [31], Khan et al. have proposed an adaptive sliding mode control approach for synchronizing complicated chaotic systems with uncertainty and disturbance. The authors of [32] have investigated an adaptive sliding control with fuzzy logic for the synchronization of chaotic fractional order systems with uncertainty and exogenous shocks. In [33], a class of complicated fractional order systems with non-uniform order has been explored, and an adaptive sliding control has been proposed for synchronizing this class of systems. An adaptive controller has been devised in [34] for integer-order synchronization with uncertainty and unknown time delays.

The time delays in the system being unknown is one of the factors to consider for uncertain systems with time delays. This issue can provide significant challenges in the controller design process for a variety of purposes, including synchronization. The amount of time delays in the system is unknown which, on the other hand, increases the complexity of the system model, which can be considered in enhancing the security level of data transmission in the field of secure telecommunications. A sliding surface based on NLFPID is proposed to direct and keep all the states of the system to the sliding surface in order to synchronize a class of chaotic systems with uncertainty and unknown time delays.

A sliding surface based on NLFPID is proposed to direct and keep all the states of the system on the sliding surface in order to synchronize a class of chaotic systems with uncertainty and unknown time delays. Then, for robust synchronization of uncertain chaotic systems with unknown time delays, a fractional order adaptive control is adopted. Following that, updating and estimating laws for uncertain parameters are determined using an appropriate Lyapunov function and Lipschitz condition to ensure the system's stability. Finally, the synchronization of two chaotic jerk systems with uncertain and unknown time delays, as well as uncertainty and distortion, are examined and simulated in order to evaluate the performance of the proposed approach. The simulation results demonstrate the efficacy of the proposed adaptive-sliding control mechanism for synchronization that is robust against uncertainty, external distortions, and unknown time delays. Furthermore, based on the obtained results, the proposed control approach is effective in estimating uncertain parameters and unknown delays.

To implement the proposed control strategy, a secure telecommunication mechanism based on chaos masking is presented at the end, indicating optimal security in sending and
high quality in retrieving information despite the presence of various uncertainties and uncertain parameters in the system structure.

The following innovations have been presented in the study that was carried out in order to synchronize two uncertain chaotic systems.

1. Using the NLFOPID sliding surface instead of conventional sliding surfaces.
2. The existence of unknown time delays.
3. The limits of uncertainty and disturbance are unknown.

Accordingly, based on the above concepts, using the appropriate Lyapunov function and update rules, a control mechanism has been proposed, which can overcome the problem of unknown time delays, uncertain uncertainty, and uncertain disturbances by properly adjusting the controller parameters.

This article is structured as follows. Section 2 presents the basic definitions used in the article. Section 3 describes the description of chaotic systems with uncertainty and unknown time delays. Section 4 presents an NLFPID-based sliding surface for directing and maintaining system states on the sliding surface. The adaptive robust control technique for the synchronization of uncertain chaotic systems with finite and uncertain time delays is detailed in Section 5 . Section 6 investigates and simulates the synchronization of two chaotic jerk systems with unknown time delays, as well as uncertainty and distortion, and Section 7 presents results based on the stated ideas.

## 2. Basic Definitions of the Fractional-Order Derivative

Definition 1. The fractional order integral and derivative are defined as follows [27]:

$$
D_{t}^{\alpha}=\left\{\begin{array}{lr}
\frac{d^{\alpha}}{d t^{\alpha}} & \alpha>0  \tag{1}\\
1 & \alpha=0 \\
\int_{a}^{t}(d \tau)^{-\alpha} & \alpha<0
\end{array}\right.
$$

where $D_{t}^{\alpha}$ is the fractional order operator.
Definition 2. The Riemann-Liouville fractional order integral order a of the function $f(t)$ is defined as follows [28]:

$$
\begin{equation*}
t_{0} I_{t}^{\alpha} f(t)=\frac{1}{\Gamma(\alpha)} \int_{t_{0}}^{t} \frac{f(\tau)}{(t-\tau)^{1-\alpha}} d \tau \tag{2}
\end{equation*}
$$

where $t_{0}$ is the initial time and $\Gamma(\alpha)$ is the Gamma function, which is defined as follows:

$$
\begin{equation*}
\Gamma(\alpha)=\int_{0}^{\infty} e^{-t} t^{\alpha-1} d t \tag{3}
\end{equation*}
$$

where is the Gamma function operator.
Definition 3. Suppose $n-1<\alpha \leq n . n \in N$. The fractional Riemann-Liouville derivative of order $\alpha$ is defined for the function $f(t)$ as follows [26]:

$$
\begin{equation*}
t_{0} D_{t}^{\alpha} f(t)=\frac{d^{\alpha} f(t)}{d t^{\alpha}}=\frac{1}{\Gamma(n-\alpha)} \frac{d^{n}}{d t^{n}} \int_{t_{0}}^{t} \frac{f(\tau)}{(t-\tau)^{\alpha-n+1}} d \tau \tag{4}
\end{equation*}
$$

Remark 1. The Riemann-Liouville fractional order derivative in Equation (4) is first integrated and then derived. Therefore, the derivative of a constant number in this definition is not equal to zero.

Definition 4. The Caputo fractional order derivative of order $\alpha$ in the continuous function $f(t)$ is defined as follows [29]:

$$
t_{0} D_{t}^{\alpha} f(t)=\left\{\begin{array}{cc}
\frac{1}{\Gamma(m-\alpha)} \int_{t_{0}}^{t} \frac{f^{(m)}(\tau)}{(t-\tau)^{\alpha-m+1}} d \tau & m-1<\alpha<m  \tag{5}\\
\frac{d^{m} f(t)}{d t^{m}} & \alpha=m
\end{array}\right.
$$

where $m$ is the first integer after $\alpha$.
Definition 5. If the function $f(x, t)$ is piecewise linear and satisfies the Lipschitz conditions, then [34]:

$$
\begin{equation*}
|f(t, x)-f(t, \hat{x})| \leq \gamma_{f}|x-\hat{x}| \quad \forall x, \hat{x} \in R^{n} \tag{6}
\end{equation*}
$$

where $f(x, t)$ is Lipschitz at $x$ and the positive constant $\gamma_{f}$, is called the Lipschitz constant.

## 3. The System Descriptor Equations

The equations characterizing a class of master-slave chaotic systems with uncertainty and unknown time delays in the presence of an unknown disturbance will be introduced in this section; the canonical form dynamics of the master system are as follows:

$$
\left\{\begin{array}{cc}
\dot{x}_{i}=x_{i+1} . & 1 \leq i \leq n-1  \tag{7}\\
\dot{x}_{n}=\sigma_{0}^{T} \mathrm{x}+f\left(x\left(t-\tau_{1}\right), t\right)+\Delta f(x(t), t)+d_{1}(t)
\end{array}\right.
$$

Equations of the slave system are as follows:

$$
\left\{\begin{array}{cc}
\dot{y}_{i}=y_{i+1} \cdot & 1 \leq i \leq n-1  \tag{8}\\
\dot{y}_{n}=\sigma_{0}^{T} \mathrm{y}+g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(y(t), t)+d_{2}(t)+u(t)
\end{array}\right.
$$

The differential equations are expressed in the forms corresponding to a number of well-known chaotic systems, such as the Van der Pol oscillator, Duffing's oscillator, Genesio-Tesi's system, Arneodo's system, etc. [35], where $x(t), y(t) \in R^{n}$ describes the dynamic states of the master and slave systems, $\sigma_{0}^{T}$ represents the constant coefficients in linear states of the system, and $f\left(x\left(t-\tau_{1}\right), t\right) \cdot g\left(y\left(t-\tau_{2}\right), t\right) \in R$ are the terms of nonlinear functions with unknown time delays with a $\tau_{1}, \tau_{2}$ delay, and $\Delta f(x(t), t), \Delta g(x(t), t)$ describes nonlinear bounded uncertainties of the master and slave systems. Additionally, $d_{1}(t), d_{2}(t)$ describes the external disturbances of the master and slave systems, and $u(t)$ is the control law applied to the slave system.

Definition 6. If the following criterion is met for the systems given in Equations (1) and (2) for all conditions influencing the system, including all initial conditions, uncertainties, and unknown time delays, as well as external disturbance, the system has robust synchronization:

$$
\begin{equation*}
\lim _{t \rightarrow \infty}\left|y_{i}(t)-x_{i}(t)\right|=\lim _{t \rightarrow \infty}\left|e_{i}(t)\right|=0 . i=1, \ldots, n \tag{9}
\end{equation*}
$$

where $e_{i}(t)$ introduces synchronization errors in the master and slave systems.
As a result, the following are the dynamic equations proposing the synchronization error for the uncertain chaotic master and slave systems with unknown time delays, as specified in Equations (1) and (2):

$$
\left\{\begin{array}{cc}
\dot{e}_{i}=e_{i+1} & 1 \leq i \leq n-1  \tag{10}\\
\dot{e}_{n}=\sigma_{0}^{T} E+g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t) \\
-f\left(x\left(t-\tau_{1}\right), t\right)-\Delta f(x(t), t)-d_{1}(t)+u(t)
\end{array}\right.
$$

where $E=\left(e_{1} \cdot e_{2} \cdots \cdot e_{n}\right)^{T}$. As a result, by initially introducing a PI sliding surface and a fractional order non-linear derivative, all states of the system should be directed to and held on the sliding surface. The system's uncertainty bounds and unknown parameters should then be estimated and updated by creating an adaptive controller. In the continuation of the robust synchronization of chaotic systems (7) and (8) in the presence of external disturbances, bounded nonlinear uncertainties, and unknown time delays, it should be performed such that the dynamics of the slave system state in a finite time conforms to the dynamic behavior. The estimation error of the unknown parameters in both chaotic systems tend to zero in any state, and the robust stability of the system in finite time is guaranteed.

Assumption 1. The unknown external disturbances $d_{1}(t), d_{2}(t)$ and unknown bounded nonlinear uncertainties $\Delta f(x(t), t)$ and $\Delta g(x(t), t)$ in the master and slave systems (7) and (8) satisfy the following conditions:

$$
\begin{align*}
|\Delta f(x(t), t)| & \leq \beta_{1} \omega_{1}(x) . \\
|\Delta g(y(t), t)| & \leq \beta_{2} \omega_{2}(y) .  \tag{11}\\
\left|d_{1}(t)\right| & \leq \rho_{1} . \\
\left|d_{2}(t)\right| & \leq \rho_{2} .
\end{align*}
$$

where $|\cdot|$ describes $l_{1}$ norm and $\beta_{2}, \beta_{1}, \rho_{2}, \rho_{1}$ are unknown real positive constants, and $\omega_{2}(\cdot), \omega_{1}(\cdot)$ are known functions.
Assumption 2. Unknown time delays introduced by non-linear functions $f\left(x\left(t-\tau_{1}\right), t\right)$, $g\left(y\left(t-\tau_{2}\right), t\right) \in R$ are represented in the general forms of (7) and (8) in the master and slave systems, for each $x(t), y(t) \in R$ and, according to (6), they satisfy the following Lipschitz condition:

$$
\begin{align*}
& \left|f\left(x\left(t-\tau_{1}\right)\right)-f\left(x\left(t-\hat{\tau}_{1}\right)\right)\right| \leq k_{1}\left|x\left(t-\tau_{1}\right)-x\left(t-\hat{\tau}_{1}\right)\right| \\
& \left|x\left(t-\tau_{1}\right)-x\left(t-\hat{\tau}_{1}\right)\right| \leq m_{1}\left|\left(t-\tau_{1}\right)-\left(\mathrm{t}-\hat{\tau}_{1}\right)\right|=m_{1}\left|\widetilde{\tau}_{1}\right| \\
\Rightarrow & \left|f\left(x\left(t-\tau_{1}\right)\right)-f\left(x\left(t-\hat{\tau}_{1}\right)\right)\right| \leq l_{1}\left|\tau_{1}-\hat{\tau}_{1}\right|=l_{1}\left|\widetilde{\tau}_{1}\right|, l_{1}=k_{1} m_{1}  \tag{12}\\
& \left|g\left(y\left(t-\tau_{2}\right)\right)-g\left(y\left(t-\hat{\tau}_{2}\right)\right)\right| \leq k_{2}\left|y\left(t-\tau_{2}\right)-y\left(t-\hat{\tau}_{2}\right)\right| \\
& \left|y\left(t-\tau_{2}\right)-y\left(t-\hat{\tau}_{2}\right)\right| \leq m_{1}\left|\left(t-\tau_{2}\right)-\left(t-\hat{\tau}_{2}\right)\right|=m_{2}\left|\widetilde{\tau}_{2}\right| \\
\Rightarrow & \left|g\left(y\left(t-\tau_{2}\right)\right)-g\left(y\left(t-\hat{\tau}_{2}\right)\right)\right| \leq l_{2}\left|\tau_{2}-\hat{\tau}_{2}\right|=l_{2}\left|\widetilde{\tau}_{2}\right|, l_{2}=k_{2} m_{2}
\end{align*}
$$

where $\tau_{1}, \tau_{2} \in R$ expresses unknown time delays, $\hat{\tau}_{1}, \hat{\tau}_{2} \in R$ estimates unknown time delays, and $l_{1}$ and $l_{2}$ are positive and uncertain constants.

## 4. The Sliding Mode Control Approach Based on NLFPID Controllers

We will provide an NLFPID-based sliding surface to synchronize the chaotic system with unknown uncertainty and time delays presented in (7) and (8). The novel NLF sliding surface is presented in accordance with the NLFPID controller structure established in [35], which enhances tracking:

$$
\begin{equation*}
s(t)=h(e) \cdot\left[k_{p} e_{n}(t)+T_{I} D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}+T_{d} D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right] . \tag{13}
\end{equation*}
$$

where $h(e)$ is a nonlinear function, defined as follows:

$$
\begin{equation*}
h(e)=k_{0}+\left(1-k_{0}\right)\|E(t)\| . \quad k_{0} \in(0,1) \tag{14}
\end{equation*}
$$

where $\|E\|$ is the first norm of the dynamic state of the system error, expressed as $\|E(t)\|=\sum_{i=1}^{n}\left|\mathrm{e}_{\mathrm{i}}\right| . \mathrm{T}_{\mathrm{I}}$ and $\mathrm{T}_{\mathrm{d}}$ are the time constants of integral and derivative phrases. The parameters $\mathrm{k}_{1 \mathrm{i}}$ and $\mathrm{k}_{2 \mathrm{i}}$ are positive constant values of the sliding surface that fulfill the intended system's stability. The following conditions must be met if the system is in sliding mode:

$$
\begin{equation*}
s(t)=0 \quad, \quad \dot{s}(t)=0 \tag{15}
\end{equation*}
$$

Therefore, the fractional order derivative of the sliding surface in Equation (13) is as follows:

$$
\begin{align*}
& \dot{s}(t)=k_{0} k_{p} \dot{e}_{n}(t) \\
&+k_{0} T_{i} D^{1-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{1+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)  \tag{16}\\
&+\left(1-k_{0}\right)\left[k_{p} \frac{d}{d t}\left(\|E(t)\| e_{n}(t)\right)\right. \\
&+T_{I} \frac{d}{d t}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right) \\
&\left.+T_{d} \frac{d}{d t}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right]=0 .
\end{align*}
$$

In this scenario, we will substitute the final dynamic of the system's integer order error specified in Equation (10) in Equation (16), yielding:

$$
\begin{align*}
& \dot{s}(t)=k_{0} k_{p}( g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t)-f\left(x\left(t-\tau_{1}\right), t\right) \\
&\left.-\Delta f(x(t), t)-d_{1}(t)+\sigma_{0}^{T} \cdot E(t)+u(t)\right) \\
&+k_{0} T_{i} D^{1-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{1+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t) \\
&+\left(1-k_{0}\right)\left[k_{p} \frac{d}{d t}\left(\|E(t)\| e_{n}(t)\right)\right. \\
&+T_{I} \frac{d}{d t}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right)  \tag{17}\\
&\left.+T_{d} \frac{d}{d t}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right]=0 .
\end{align*}
$$

The control signal in this scenario is as follows:

$$
\begin{align*}
& u(t)=\frac{-1}{k_{0} k_{p}}\left(k_{0} T_{i} D^{1-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{1+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right. \\
&+\left(1-k_{0}\right)\left[k_{p} \frac{d}{d t}\left(\|E(t)\| e_{n}(t)\right)\right. \\
&+T_{I} \frac{d}{d t}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right)  \tag{18}\\
&\left.\left.+T_{d} \frac{d}{d t}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right]\right) \\
&+f\left(x\left(t-\hat{\tau}_{1}\right), t\right)-g\left(y\left(t-\hat{\tau}_{2}\right), t\right)-\sigma_{0}^{T} \cdot E(t)-b s+\bar{u}(t)
\end{align*}
$$

In Equation (18), the phrase $\bar{u}(t)$ comprises the sentences arising from estimating the bounds of the system's uncertainties and disturbances, which are specified as follows:

$$
\begin{equation*}
\left.\bar{u}(t)=-\operatorname{sgn}(s)\left[\hat{\beta}_{2} \omega_{2}(\mathrm{y})+\hat{\beta}_{1} \omega_{1}(\mathrm{x})+\hat{\rho}_{2}+\hat{\rho}_{1}\right)\right] . \tag{19}
\end{equation*}
$$

## 5. Stability Analysis of the Proposed Mechanism

In this section, the robust adaptive controller is designed using the sliding surface based on NLFPID such that the synchronization of chaotic systems is guaranteed by the proposed control approach.

Theorem 1. Synchronization of systems (6) and (7) is guaranteed by the definition of the controller $u(t)$ in spite of disturbances $d_{1}$ and $d_{2}$ and uncertainties $\Delta f$ and $\Delta g$ along with unknown time delays $\tau_{1}$ and $\tau_{2}$ :

$$
\begin{align*}
u(t)=-g(y(t- & \left.\left.\hat{\tau}_{1}\right)\right)+f\left(x\left(t-\hat{\tau}_{2}\right)\right) \\
& -\frac{1}{k_{0} k_{p}}\left(k_{0} T_{l} D^{1-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{1+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right. \\
& +\left(1-k_{0}\right)\left[k_{p} \frac{d}{d t}\left(\|E(t)\| e_{n}(t)\right)\right.  \tag{20}\\
& +T_{I} \frac{d}{d t}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right) \\
& \left.\left.+T_{d} \frac{d}{d t}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right]\right)-\sigma_{0}^{T} \cdot E(t)-b s \\
& -\operatorname{sgn}(s)\left(\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right) .
\end{align*}
$$

where $\hat{\rho}_{1}$ and $\hat{\rho}_{2}$ are the estimations of input disturbances, $\hat{\tau}_{1}$ and $\hat{\tau}_{2}$ are estimates of time delays, and $\hat{\beta}_{1} \omega_{1}$ and $\hat{\beta}_{2} \omega_{2}$ are estimates of uncertainty in the master and slave systems. Therefore, in order to guarantee the stability of the system, we use the update laws to estimate the mentioned parameters as follows:

$$
\begin{gather*}
\widetilde{\tau}_{i}=\tau_{i}-\hat{\tau}_{i} \cdot \widetilde{\beta}_{i}=\beta_{i}-\hat{\beta}_{i} \cdot \widetilde{\rho}_{i}=\rho_{i}-\hat{\rho}_{i} . \\
\dot{\hat{\tau}}_{i}=-\dot{\widetilde{\tau}}_{i}=|s| \operatorname{sgn}\left(\widetilde{\tau}_{i}\right)+b \widetilde{\tau}_{i} . \\
\dot{\hat{\rho}}_{i}=-\dot{\tilde{\rho}}_{i}=k_{0} k_{p}|s|+b \widetilde{\rho}_{i} .  \tag{21}\\
\dot{\hat{\beta}}_{2}=-\dot{\widetilde{\beta}}_{2}=k_{0} k_{p}|s| \omega_{2}(y)+b \widetilde{\beta}_{2} . \\
\dot{\hat{\beta}}_{1}=-\dot{\widetilde{\beta}}_{1}=k_{0} k_{p}|s| \omega_{1}(x)+b \widetilde{\beta}_{1} .
\end{gather*}
$$

Proof. The following Lyapunov function is defined as follows:

$$
\begin{equation*}
v(t)=\frac{1}{2}\left[s^{2}(t)+\widetilde{\beta}_{1}^{2}+\widetilde{\beta}_{2}^{2}+l_{1} \widetilde{\tau}_{1}^{2}+l_{2} \widetilde{\tau}_{2}^{2}+\widetilde{\rho}_{1}^{2}+\widetilde{\rho}_{2}^{2}\right] \tag{22}
\end{equation*}
$$

According to Equation (22), the derivative of the Lyapunov function is as follows:

$$
\begin{align*}
\Rightarrow \dot{v}(t)= & \frac{1}{2} \frac{d}{d t}\left(s^{2}+\widetilde{\beta}_{1}^{2}+\widetilde{\beta}_{2}^{2}+l_{1} \widetilde{\tau}_{1}^{2}+l_{2} \widetilde{\tau}_{2}^{2}+\widetilde{\rho}_{1}^{2}+\widetilde{\rho}_{2}^{2}\right) \\
& =s \dot{s}+\sum_{i=1}^{2}\left(\widetilde{\beta}_{i} \dot{\tilde{\beta}}_{i}+l_{i} \widetilde{\tau}_{i} \dot{\tilde{\tau}}_{i}+\widetilde{\rho}_{i} \dot{\tilde{\rho}}_{i}\right) \tag{23}
\end{align*}
$$

By substituting (17) in (23), Equation (24) is obtained:

$$
\begin{align*}
\dot{v}(t)=s \cdot\left[k_{0} k_{p}\right. & \left(g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t)\right. \\
& \left.-\left(f\left(x\left(t-\tau_{1}\right), t\right)+\Delta f(x(t), t)+d_{1}(t)\right)+\sigma_{0}^{T} \cdot E(t)+u(t)\right) \\
& +k_{0} T_{i} D^{1-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{1+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t) \\
+ & \left(1-k_{0}\right)\left[k_{p} \frac{d}{d t}\left(\|E(t)\| e_{n}(t)\right)\right.  \tag{24}\\
& +T_{I} \frac{d}{d t}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right) \\
& \left.\left.+T_{d} \frac{d}{d t}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right]\right]+\sum_{i=1}^{2}\left(\widetilde{\beta}_{i} \dot{\widetilde{\beta}}_{i}+l_{i} \widetilde{\tau}_{i} \dot{\widetilde{\tau}}_{i}+\widetilde{\rho}_{i} \dot{\tilde{\rho}}_{i}\right)
\end{align*}
$$

By substituting (20) in (24), Equation (25) is obtained:

$$
\begin{align*}
\dot{v}(t)=s \cdot k_{0} k_{p}(g( & \left.y\left(t-\tau_{2}\right), t\right)-g\left(y\left(t-\hat{\tau}_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t) \\
& +f\left(x\left(t-\hat{\tau}_{1}\right), t\right)-f\left(x\left(t-\tau_{1}\right), t\right)-\Delta f(x(t), t)-d_{1}(t)-b s \\
& \left.-\operatorname{sgn}(s)\left[\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right]\right)  \tag{25}\\
& +\sum_{i=1}^{2}\left(\widetilde{\beta}_{i} \dot{\tilde{\beta}}_{i}+l_{i} \widetilde{\tau}_{i} \dot{\widetilde{\tau}}_{i}+\widetilde{\rho}_{i} \dot{\tilde{\rho}}_{i}\right)
\end{align*}
$$

Thus, we have:

$$
\begin{align*}
& \dot{v}(t) \leq|s| \cdot\left[k _ { 0 } k _ { p } \left(\left|g\left(y\left(t-\tau_{2}\right), t\right)-g\left(y\left(t-\hat{\tau}_{2}\right), t\right)\right|+|\Delta g(x(t), t)|\right.\right. \\
&+\left|f\left(x\left(t-\hat{\tau}_{1}\right), t\right)-f\left(x\left(t-\tau_{1}\right), t\right)\right|-|\Delta f(x(t), t)| \\
&\left.\left.+\left|d_{2}(t)-d_{1}(t)\right|\right)\right]-k_{0} k_{p} b s^{2}  \tag{26}\\
&\left.+k_{0} k_{p} s\left(-\operatorname{sgn}(s)\left[\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right)\right]\right) \\
&+\sum_{i=1}^{2}\left(\widetilde{\beta}_{i} \dot{\widetilde{\beta}}_{i}+l_{i} \widetilde{\tau}_{i} \dot{\tilde{\tau}}_{i}+\widetilde{\rho}_{i} \dot{\tilde{\rho}}_{i}\right)
\end{align*}
$$

By substituting (11) and (12) in (26), Equation (27) is obtained:

$$
\begin{align*}
\dot{v}(t) \leq|s| \cdot & {\left[k_{0} k_{p}\left(l_{2}\left|\tau_{2}-\hat{\tau}_{2}\right|+\beta_{2} \omega_{2}(y)+l_{1}\left|\tau_{1}-\hat{\tau}_{1}\right|+\beta_{1} \omega_{1}(x)+\rho_{1}-\rho_{2}\right)\right] } \\
& -k_{0} k_{p} b s^{2} \\
& \left.+k_{0} k_{p} s\left(\left(-\operatorname{sgn}(s)\left[\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right)\right]\right)\right)  \tag{27}\\
& +\sum_{i=1}^{2}\left(\widetilde{\beta}_{i} \dot{\tilde{\beta}}_{i}+l_{i} \widetilde{\tau}_{i} \dot{\tilde{\tau}}_{i}+\widetilde{\rho}_{i} \dot{\tilde{\rho}}_{i}\right) .
\end{align*}
$$

The derivative of the Lyapunov function would be as follows:

$$
\begin{align*}
\dot{v}(t) \leq|s|\left[k _ { 0 } k _ { p } \left(l_{1}\left|\widetilde{\tau}_{1}\right|\right.\right. & \left.\left.+\widetilde{\beta}_{2} \omega_{2}(y)+l_{2}\left|\widetilde{\tau}_{2}\right|+\widetilde{\beta}_{1} \omega_{1}(y)+\hat{\rho}_{2}+\hat{\rho}_{1}\right)\right]-b s^{2} \\
& +\sum_{i=1}^{2}\left(\widetilde{\beta}_{i} \dot{\tilde{\beta}}_{i}+l_{i} \widetilde{\tau}_{i} \dot{\tilde{\tau}}_{i}+\widetilde{\rho}_{i} \dot{\tilde{\rho}}_{i}\right) . \tag{28}
\end{align*}
$$

Using the update laws (21) and substituting them in (28), Equation (29) is obtained:

$$
\begin{equation*}
\Rightarrow \dot{v}(t) \leq-b\left(s^{2}+\widetilde{\beta}_{1}^{2}+\widetilde{\beta}_{2}^{2}+\widetilde{\tau}_{1}^{2}+\widetilde{\tau}_{2}^{2}+\widetilde{\rho}_{1}^{2}+\widetilde{\rho}_{2}^{2}\right)=-b v(t) . \tag{29}
\end{equation*}
$$

Thus, the stability of the proposed method for synchronization of the chaotic system with uncertainty, disturbance, and unknown time delays is proved.

## 6. Simulation Results

This section evaluates the accuracy of synchronizing uncertain chaotic systems with unknown time delays using the proposed control mechanism based on an NLF sliding surface, adaptive controller, and update rules that estimate the system's parameters. Two modified jerk chaotic systems with the aforementioned properties were used for this purpose. The equations governing the master system in the canonical form are as follows [33]:

$$
\left\{\begin{array}{c}
\dot{x}_{1}=x_{2}  \tag{30}\\
\dot{x}_{2}=x_{3} \\
\dot{x}_{3}=-\varepsilon_{1} x_{1}(t)-x_{2}(t)-\varepsilon_{2} x_{3}(t)+f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right)
\end{array}\right.
$$

where $\varepsilon_{1}=\frac{3}{2}, \varepsilon_{2}=0.35$, and $f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right)$ is a piecewise linear function:

$$
\begin{align*}
f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right) & =\frac{1}{2}\left(v_{0}-v_{1}\right)\left[\left|x_{1}\left(t-\tau_{1}\right)+1\right|-\left|x_{1}\left(t-\tau_{1}\right)-1\right|\right]  \tag{31}\\
& +v_{1} x_{1}\left(t-\tau_{1}\right) .
\end{align*}
$$

where $v_{0}<-1<v_{1}<0$ and $v_{0}=-2.5, v_{1}=-0.5$.
When the initial conditions are selected as $\left(x_{1}(0) ; x_{2}(0) ; x_{3}(0)\right)^{T}=(-0.52 ; 0.52 ; 0.87)^{T}$, the system's chaotic behavior would be as shown in Figure 1.

$$
\left\{\begin{array}{c}
\dot{x}_{1}=x_{2}  \tag{32}\\
\dot{x}_{2}=x_{3} \\
\dot{x}_{3}=-\varepsilon_{1} x_{1}(t)-x_{2}(t)-\varepsilon_{2} x_{3}(t)+f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right)+\Delta f(x(t), t)+d_{1}(t)
\end{array}\right.
$$

The dynamic of the slave system follows the equations below:

$$
\left\{\begin{array}{c}
\dot{y}_{1}=y_{2}  \tag{33}\\
\dot{y}_{2}=y_{3} \\
\dot{y}_{3}=-\varepsilon_{1} y_{1}(t)-y_{2}(t)-\varepsilon_{2} y_{3}(t)+g_{3}\left(y_{1}\left(t-\tau_{2}\right), t\right)+\Delta g(y(t), t)+d_{2}(t)+u(t)
\end{array}\right.
$$

where the nonlinear terms of the slave system are as follows:

$$
\begin{gather*}
g_{3}\left(y_{1}\left(t-\tau_{2}\right), t\right)=\frac{1}{2}\left(v_{0}-v_{1}\right)\left[\left|y_{1}\left(t-\tau_{2}\right)+1\right|-\left|y_{1}\left(t-\tau_{2}\right)-1\right|\right]  \tag{34}\\
+v_{1} y_{1}\left(t-\tau_{2}\right)
\end{gather*}
$$

Accordingly, the dynamic of the synchronization error for the chaotic jerk master and slave systems would be as follows:

$$
\left\{\begin{array}{c}
\dot{e}_{1}=e_{2}  \tag{35}\\
\dot{e}_{2}=e_{3} \\
\dot{e}_{3}=-\varepsilon_{1} e_{1}(t)-e_{2}(t)-\varepsilon_{2} e_{3}(t)-g\left(y_{1}\left(t-\tau_{2}\right)\right)+f\left(x_{1}\left(t-\tau_{1}\right)\right) \\
+\Delta g(y(t), t)-\Delta f(x(t), t)+d_{2}(t)-d_{1}(t)+u(t)
\end{array}\right.
$$

In this step, we apply the robust adaptive control signal, which is designed by combining the sliding surface based on the NLFPID controllers and described in Equation (20), to the slave system.

The following figures show the behavior of the chaotic system synchronized with the above dynamic equations before and after applying the proposed control signal. Figures 1 and 2 respectively show the phase diagram and the behavior of the jerk system without applying the controller.


Figure 1. Phase diagram of the jerk master and slave systems without applying the controller.


Figure 2. Behavior of master and slave system states without applying the control signal.
Figure 3 shows the behavior of the synchronized system. The synchronization error of the jerk system is shown in Figure 4. The control signal based on the proposed mechanism
is shown in Figure 5. The estimation of the system parameters is shown in Figure 6. The disturbance and uncertainty of the master and slave system are shown in Figure 7.


Figure 3. Synchronization of jerk systems using the proposed mechanism and applying the control signal at $\mathrm{t}=5 \mathrm{~s}$.


Figure 4. Synchronization error of the master and slave systems.


Figure 5. The control signal based on the proposed adaptive sliding mode control.


Figure 6. Estimation error of the system parameters.

Disturbances and Uncertaintis to Master-Slave Systems


Figure 7. Uncertainty and disturbances of the master and slave systems.
In this article, simulations have been performed for $t=40 \mathrm{~S}$, where $k_{11}=k_{22}=10$ and $k_{12}=k_{21}=20$ have been selected. Additionally, the gain and time constants of the non-linear fractional order PID sliding surface are $k_{p}=1.5, T_{i}=0.75$ and $T_{d}=0.5$. The fractional order of the integral and derivative part of the sliding surface is defined as $\delta=0.4$ and $\lambda=0.75$. The parameters of the proposed robust controller are $b=2$. The uncertain time delays of the system are $\tau_{1}=0.65$ and $\tau_{2}=0.35$. Unknown disturbances applied to both systems are as follows:

$$
d_{1}(t)=\sin ^{2} 3 t+2 \cos 4 t, \quad d_{2}(t)=\sin 2 t+0.4 \sin \pi t
$$

## 7. Application in Secure Communication

The synchronization of the chaotic system with uncertainty and unknown time delays was discussed in the preceding section. The proposed robust control strategy is placed in the framework of a secure communication mechanism in this part such that the message signal is sent by the master system after the encryption process. The chaotic masking method was utilized for this aim. Figure 8 depicts how to convey the primary message signal using the proposed mechanism. A wireless communication channel is employed in this block diagram. The communication channel might be either wired or wireless.

If $M(t)$ is the original message signal that is coupled with the master system and $S(t)$ is the sent message, then [36]:

$$
\begin{equation*}
S(t)=M(t)+\sum_{j=1}^{n} \gamma_{j} x_{j} . \tag{36}
\end{equation*}
$$

where $x_{j}$ is the state of the chaotic system and $\gamma_{j}$ is the constant number. $S(t)$ is the masked chaotic signal that is sent by the created communication channel. Using the proposed chaotic synchronization mechanism, the system states are integrated with the message signal in a weighted manner to perform the masking process. The receiver side can receive the message signal as follows:

$$
\begin{equation*}
R(t)=S(t)-\sum_{j=1}^{n} \gamma_{j} y_{j} \tag{37}
\end{equation*}
$$



Figure 8. Chaotic secure communication structure based on the proposed approach.
Using the concept of synchronization, the following signal can be reconstructed in the receiver as follows:

$$
\begin{gather*}
R(t)=M(t)+\sum_{j=1}^{n} \gamma_{j} x_{j}-\sum_{j=1}^{n} \gamma_{j} y_{j}=M(t)+\sum_{j=1}^{n} \gamma_{j}\left(x_{j}-y_{j}\right)  \tag{38}\\
\Rightarrow x_{j}-y_{j}=e_{j} \cong 0 \Rightarrow R(t) \rightarrow M(t) .
\end{gather*}
$$

In this step, two message signals are generated, which are applied to the synchronized system via the chaotic masking mechanism mentioned above. In this phase, the control signal is applied to the slave system, and the message signal is applied to the synchronized system after achieving the established time to reduce the mean square error. The transition time in the simulations ends at $t=4.5 \mathrm{~s}$. At this point, the encrypted message's signal is applied, and the signal $S(t)$ is received on the receiver's side and decoded to extract the signal $R(t)$. The simulation results in the Figures 9 and 10 demonstrate a good performance of the proposed system.

Table 1 shows the mean square error of the original message signal and the recovered signal.


Figure 9. Message 1, masked and retrieved, and the error between these two signals.


Figure 10. Message 2, masked and retrieved, and the error between these two signals.
Table 1. Represents the mean square error of the main message signal and the retrieved signal.

|  | Root Mean Squre Error |  |
| :---: | :---: | :---: |
| $M_{1}$ and $R_{1}$ |  | $M_{2}$ and $R_{2}$ |
| 0.041544 | 0.025857 |  |

The following are the message signals:

$$
\begin{gathered}
M_{1}(t)=1.8 \sin (1.7 t)+4.5 \cos (10.2 t)+5.4 \sin (1.33 \pi t) \\
M_{2}(t)=2.25 \sin (1.9 t)+2.88 \cos (10.66 t)+4.05 \sin (2.33 \pi t)
\end{gathered}
$$

The masked parameters of the message signal are $\gamma_{1}=1, \gamma_{2}=0.5$, and $\gamma_{3}=0.33$. Thus, encoding is as follows:

$$
\begin{gathered}
S_{1}(t)=M_{1}+\gamma_{1} x_{1} \\
S_{2}(t)=M_{2}+\gamma_{2} x_{1}+\gamma_{3} x_{2}
\end{gathered}
$$

## 8. Conclusions

This study investigates the robust synchronization of a class of chaotic systems with uncertainty, external disturbances, and unknown parameters, such as unknown time delays, by introducing a new adaptive sliding mode control technique. First, an NLFPID-based sliding surface is proposed in the suggested robust control mechanism. The adaptive laws are then defined in order to estimate the uncertain parameters of the system using Lyapunov theory and Lipschitz conditions in chaotic systems, and ultimately the stability of the proposed robust control system is proven. The synchronization of two uncertain jerk chaotic systems with unknown time delays based on the proposed controller is simulated using MATLAB, and the results express the capability and desired performance of the proposed approach in the robust synchronization of the mentioned systems. Finally, the proposed adaptive sliding mode control approach has been used in a robust secure chaotic communication mechanism, and the simulation results indicate favorable quality in the secure transmission and reception of information despite uncertain parameters in the master and slave systems of the communication mechanism.

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Article

# Robust Multi-Criteria Traffic Network Equilibrium Problems with Path Capacity Constraints 

Xing-Xing Ma and Yang-Dong Xu *<br>Department of Mathematics, Chongqing University of Posts and Telecommunications, Chongqing 400065, China; maxx0205@126.com<br>* Correspondence: xyd04010241@126.com

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#### Abstract

With the progress of society and the diversification of transportation modes, people are faced with more and more complicated travel choices, and thus, multi-criteria route choosing optimization problems have drawn increased attention in recent years. A number of multi-criteria traffic network equilibrium problems have been proposed, but most of them do not involve data uncertainty nor computational methods. This paper focuses on the methods for solving robust multi-criteria traffic network equilibrium problems with path capacity constraints. The concepts of the robust vector equilibrium and the robust vector equilibrium with respect to the worst case are introduced, respectively. For the robust vector equilibrium, an equivalent min-max optimization problem is constructed. A direct search algorithm, in which the step size without derivatives and redundant parameters, is proposed for solving this min-max problem. In addition, we construct a smoothing optimization problem based on a variant version of ReLU activation function to compute the robust weak vector equilibrium flows with respect to the worst case and then find robust vector equilibrium flows with respect to the worst case by using the heaviside step function. Finally, extensive numerical examples are given to illustrate the excellence of our algorithms compared with existing algorithms. It is shown that the proposed min-max algorithm may take less time to find the robust vector equilibrium flows and the smoothing method can more effectively generate a subset of the robust vector equilibrium with respect to the worst case.


Keywords: multi-criteria traffic network; robust vector equilibrium; min-max method; smoothing method

## 1. Introduction

Traveling is necessary for everyday human life. However, with the progress of society and the diversification of transportation modes, people also expect to find the most efficient route. Traffic network equilibrium problem can describe the distributions of traffic flows in the logistics industry and transportation network, which is expected to provide an effective method for travelers to choose an optimization route. The fundamental principle in the model is the concept of equilibrium that was initially introduced by Wardrop [1]. The principle asserts that travelers will choose the path only if the cost for this path is the minimum possible among all the paths joining the same O-D pair.

### 1.1. Literature Reviews

It has been shown that the Wardrop equilibrium concept is a powerful principle which is widely used in supply and demand networks, traffic assignment, optimization of traffic control, and other fields (see, e.g., Athanasenas [2]; Nagurney [3]; Ji and Chu [4]; Xu et al. [5]; Wang et al. [6]; Ma et al. [7]). It is worth noting that most of these equilibrium models in the above references are based on a single criterion. Travelers (in this paper, we use the terms 'user' and 'traveler' interchangeably) will naturally consider multi-criteria when choosing travel paths, including travel time, distance, cost, weather, safety, and other relevant factors. The equilibrium model with multi-criteria was first put forward by Chen and Yen [8],
which was an extension of the classical Wardrop user equilibrium principle. Regarding the theoretical analysis for multi-criteria traffic equilibrium models, we refer the reader to Yang and Goh [9], Li et al. [10], Luc et al. [11], and Raith and Ehrgott [12].

Recently, Phuong and Luc [13] established the equivalent relationship between strong vector equilibrium flows and the solutions of variational inequality problems in terms of a kind of increasing functions. Moreover, they presented a modified projection method to handle multi-criteria traffic network equilibrium problems. Subsequently, Luc and Phuong [14] introduced two optimization problems to show that the optimal solutions are exactly the equilibrium of the traffic network and then put forward a modified FrankWolfe gradient algorithm for multi-criteria traffic network equilibrium problems. However, this method may lead to the non-differentiability of the objective functions of the two optimization problems. After that, Phuong [15] proposed a smoothing method to solve multi-criteria network equilibrium problems. Although this method solves the defects existing in [14], it does not take into account the data uncertainty.

In the actual traffic network, there are various uncertain factors, such as traveler preferences, weather, traffic congestion, and holidays. Hence, uncertainty in the logistics industry and transportation has received more and more attention. In recent years, some related works with uncertain demands or uncertain parameters in traffic network or ecological networks have been investigated in [16-21]. Daniele and Giuffré [16] investigated a general random traffic equilibrium problem and characterize the random Wardrop equilibrium distribution by means of a random variational inequality. Dragicevic and Gurtoo [17] modeled the maintenance of ecological networks in forest environments based on random processes, such as extreme natural events. However, the two above papers do not consider the multi-criteria. Recently, Ehrgott and Wang [18] presented alternative approaches for combining the principles of multi-objective decision-making with a stochastic user equilibrium model based on random utility theory. However, since uncertain parameters in [18] need to know probabilistic information, this may be inconsistent with the reality because the probabilistic information of related data are usually known. Cao et al. [19] and Wei et al. [20] only discussed relationships between and the solutions of variational inequality and robust equilibrium flows but not give the computational methods. Minh and Phuong [21] paid attention to a modified Frank-Wolfe gradient algorithm for robust equilibrium flows. The uncertain data in the model proposed in [21] are in a parameter set that does not need probability information. However, the computational efficiency of the algorithm is not good, due to the non-differentiability of the objective functions. In all, there are some research gaps on computational methods for robust multi-criteria traffic network equilibrium problems with path capacity constraint. This prompts us to continuously investigate this topic.

### 1.2. Contributions

To overcome computational inefficiency for the robust vector equilibrium flows in existing methods, this paper proposes two new computational methods for the robust vector equilibrium principle and the robust vector equilibrium principle with respect to the worst case, respectively. Firstly, an equivalent min-max optimization problem is constructed, in which the solution is equivalent to the robust vector equilibrium flow. A direct search algorithm with constraints for solving this problem is proposed. For the robust vector equilibrium with respect to the worst case, we transform it into an deterministic vector equilibrium problem based on a variant version of ReLU activation function. Then, we give an algorithm to solve the robust vector equilibrium with respect to the worst case.

In summary, the contributions of the manuscript are ranked in ascending gathered as follows:
(1) The robust vector equilibrium and the robust vector equilibrium with respect to the worst case principles are introduced.
(2) An equivalent min-max optimization problem is established and then a direct search algorithm is proposed to generate a subset of robust vector equilibrium flows.
(3) To generate a subset of the robust vector equilibrium with respect to the worst case, a two-step strategy is implemented. More specifically, a smoothing optimization problem is constructed based on a variant version of ReLU activation function to compute the robust weak vector equilibrium flows with respect to the worst case, and then, the robust vector equilibrium flows are found with respect to the worst case by using the heaviside step function.
This paper is divided into the following parts. Section 2 mainly introduces the robust vector equilibrium principle and robust vector equilibrium principle with respect to the worst case. Section 3 gives a min-max method to generate the subset of robust vector equilibrium flows. Section 4 presents a smoothing algorithm to find the subset of the robust vector equilibrium principle with respect to the worst case. Finally, conclusions of this paper and discussions for future research are provided in Section 5.

## 2. Definitions and Main Derivations

We review some fundamental definitions and properties that are relevant to this study. Throughout this paper, let $\mathbb{R}^{*}(*=n, m)$ denote the $*$-dimensional Euclidean space. Let $\mathbb{R}_{+}^{m}:=\left\{x \in \mathbb{R}^{m}: x_{i} \geq 0, i=1, \cdots, m\right\}$ and $\mathbb{R}_{++}^{m}:=\left\{x \in \mathbb{R}^{m}: x_{i}>0, i=1, \cdots, m\right\}$. The superscript $\top$ denotes transpose. The partial order in $\mathbb{R}^{m}$ is induced by $\mathbb{R}_{+}^{m}$, defined by:

$$
x \geqq y \text { if } x_{i} \geq y_{i} \text { for all } i=1, \ldots, m
$$

$x \succeq y$ if $x_{i} \geq y_{i}$ for all $i=1, \ldots, m$ and there exists $i_{0}$ such that $x_{i_{0}}>y_{i_{0}}$. and the following stronger relation is given by:

$$
x \succ y \text { if } x_{i}>y_{i} \text { for all } i=1, \ldots, m
$$

Next, we will denote by $e$ the vector of all ones. Given $X \subseteq \mathbb{R}^{m}$, the set of minimal elements of $X$ is denoted by $\operatorname{Min}(X)$, consists of vectors $x \in X$ such that there is no $x^{\prime} \in X, x^{\prime} \preceq x$.

Definition 1. Given $f: \mathbb{R} \times \mathbb{R} \rightarrow \mathbb{R}$, we say that a point $\left(x^{*}, y^{*}\right)$ is a saddle-point of the function $f$, if

$$
f\left(x^{*}, y\right) \leq f\left(x^{*}, y^{*}\right) \leq f\left(x, y^{*}\right), \forall(x, y) \in \mathbb{R} \times \mathbb{R} .
$$

### 2.1. Robust Multi-Criteria Traffic Network

For a traffic network, $\mathcal{N}$ denotes the set of the nodes and $\mathcal{E}$ denotes the set of directed arcs. Let $\mathcal{W}$ be the set of origin-destination (O-D) pairs and $\mathcal{D}=\left(d_{\omega}\right)_{\omega \in \mathcal{W}}$ be the demand vector, where $d_{\omega}>0$ is the flow demand on O-D pair $\omega$. Thus, a traffic network is always denoted by $G=\{\mathcal{N}, \mathcal{E}, \mathcal{W}, \mathcal{D}\}$. For $\omega \in \mathcal{W}, P_{\omega}$ is the set of available paths on the O-D pair $\omega$ and $P=\cup_{\omega \in \mathcal{W}} P_{\omega}$ is the set of all available paths of the network. Let $n=\sum_{\omega \in \mathcal{W}}\left|p_{\omega}\right|$. For a given $p_{k} \in P_{\omega}, y_{p_{k}}$ is the traffic flow on this path and $y=\left(y_{1}, y_{2}, \cdots, y_{n}\right)^{\top} \in \mathbb{R}^{n}$ is called a path flow. For given $p_{k} \in P_{\omega}$, suppose $l_{p_{k}} \in \mathbb{R}_{+}, u_{p_{k}} \in \mathbb{R}_{+}$with $l_{p_{k}}<u_{p_{k}}$; the path flow $y_{p_{k}}$ needs to satisfy the capacity constraint $l_{p_{k}} \leq y_{p_{k}} \leq u_{p_{k}}$. The traffic load is always presented by arc flows $z_{\alpha}, \alpha \in \mathcal{E}$, or path flows $y_{p_{k}}, p_{k} \in P$. Given a path flow, the arc flow can be obtained by the following formula:

$$
z_{\alpha}=\sum_{p_{k} \in P} y_{p_{k}} \delta_{\alpha p_{k}}
$$

where

$$
\delta_{\alpha p_{k}}= \begin{cases}1, & \text { if } \alpha \text { belongs to path } p_{k} \\ 0, & \text { otherwise }\end{cases}
$$

The arc flow is denoted by $z:=\left(z_{\alpha}\right)_{\alpha \in \mathcal{E}}$. A path flow $y$ is said to be feasible flow if it satisfies:

$$
\Omega=\left\{y \in \mathbb{R}_{+}: \forall \omega \in \mathcal{W}, \forall p_{k} \in P_{\omega}, l_{p_{k}} \leq y_{p_{k}} \leq u_{p_{k}} \sum_{p_{k} \in P_{\omega}} y_{p_{k}}=d_{\omega}\right\}
$$

Let $t_{\alpha}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ be a vector-valued cost function along with arc $\alpha \in \mathcal{E}$. Let $c_{p_{k}}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ be a vector-valued cost function on the path $p_{k}$. Thus, we have that the cost function $c_{p_{k}}$ for path $p_{k}$ is the sum of cost functions for arcs belonging to path $p_{k}$, namely:

$$
\begin{equation*}
c_{p_{k}}(y)=\sum_{\alpha \in \mathcal{E}} \delta_{\alpha p_{k}} t_{\alpha}(y) \tag{1}
\end{equation*}
$$

However, the path cost functions may be perturbed in reality. This means that it not only depends on the path flow $y$ but also on parameters of $\xi \in U:=U_{1} \times U_{2} \times \cdots \times U_{n}$. Throughout this paper, the cost function $c_{p_{k}}(y, \xi)$ is often given in the form $c_{p_{k}}(y, \xi)=c_{p_{k}}(y)+\xi p_{k}$.

### 2.2. Robust Vector Equilibrium and Related Concepts

Now, we give the following definitions on a robust vector equilibrium and a robust (weak) vector equilibrium with respect to the worst case.

Definition 2. A feasible flow $\bar{y} \in \Omega$ is said to be a robust vector equilibrium, if for each $O-D$ $\omega \in W$, path $p_{k}, p_{j} \in P_{\omega}$, one has:

$$
c_{p_{k}}(\bar{y}, \xi)-c_{p_{j}}(\bar{y}, \xi) \succeq 0_{\mathbb{R}^{m}}, \forall \xi \in U \Rightarrow \text { either } \bar{y}_{p_{k}}=l_{p_{k}} \text { or } \bar{y}_{p_{j}}=u_{p_{j}}
$$

The worst case of the cost function on the path $p_{k}$ under all possible scenarios is defined as follows:

$$
C_{p_{k}}(y)=\left(\begin{array}{c}
\sup _{\xi \in U} c_{1 p_{k}}(y, \xi) \\
\vdots \\
\sup _{\xi \in U} c_{m p_{k}}(y, \xi)
\end{array}\right), C_{p_{j}}(y)=\left(\begin{array}{c}
\sup _{\xi \in U} c_{1 p_{j}}(y, \xi) \\
\vdots \\
\sup _{\tilde{\xi} \in U} c_{m p_{j}}(y, \xi)
\end{array}\right)
$$

The following definitions are given based on the worst case of path costs, which is called the robust vector equilibrium with respect to the worst case and the robust weak vector equilibrium with respect to the worst case.

Definition 3. A feasible flow $\bar{y} \in \Omega$ is a robust vector equilibrium with respect to the worst case, if for $\forall \omega \in W, \forall p_{k}, p_{j} \in P_{\omega}$, one has:

$$
C_{p_{k}}(\bar{y})-C_{p_{j}}(\bar{y}) \succeq 0_{\mathbb{R}^{m}} \Rightarrow \text { either } \bar{y}_{p_{k}}=l_{p_{k}} \text { or } \bar{y}_{p_{j}}=u_{p_{j}}
$$

Definition 4. A feasible flow $\bar{y} \in \Omega$ is in robust weak vector equilibrium with respect to the worst case, if for $\omega \in W, p_{k}, p_{j} \in P_{\omega}$, one has:

$$
C_{p_{k}}(\bar{y})-C_{p_{j}}(\bar{y}) \succ 0_{\mathbb{R}^{m}} \Rightarrow \text { either } \bar{y}_{p_{k}}=l_{p_{k}} \text { or } \bar{y}_{p_{j}}=u_{p_{j}} .
$$

Remark 1. What should be noteworthy is that a robust vector equilibrium with respect to the worst case is also a robust vector equilibrium when $U$ is a compact set. Conversely, it is not necessarily true. Although there is no parameter in the concept of the robust vector equilibrium with respect to the worst case, it still depends on the values of the parameter or sensitive to parameters. Now, we give the following example to illustrate the above cases.

Example 1. Consider a network problem with one $O-D$ pair $\omega=\left(x, x^{\prime}\right)$. Two criteria, i.e., travel time and travel cost, and two available paths, i.e., $P_{\omega}=\left\{p_{1}, p_{2}\right\}$, with the travel demand $d_{\omega}=30$.

Assume that the path capacity constraints and cost function on the paths $p_{1}$ and $p_{2}$ are, respectively, given as follows:

$$
\begin{gathered}
l_{p_{1}}=0, l_{p_{2}}=0 ; u_{p_{1}}=30, u_{p_{2}}=30 . \\
c_{p_{1}}\left(y, \xi_{1}\right)=\binom{y_{1}+2 y_{2}+\xi_{1}}{6 y_{1}+2 y_{2}+\xi_{1}}, \quad c_{p_{2}}\left(y, \xi_{2}\right)=\binom{y_{1}+6 y_{2}}{6 y_{1}+2 y_{2}-\xi_{2}} .
\end{gathered}
$$

with $\xi_{1} \in[-1,2]$ and $\xi_{2} \in[0,1]$.
Direct computation shows that $\bar{y}=(30,0)$ is the robust vector equilibrium. However, it is not the robust vector equilibrium with respect to the worst case since we have:

$$
C_{p_{1}}(\bar{y})=\binom{32}{182}, \quad C_{p_{2}}(\bar{y})=\binom{30}{180}
$$

but $y_{p_{1}} \neq l_{p_{1}}$ and $y_{p_{2}} \neq u_{p_{2}}$.
If $\xi_{1} \in[-1,0]$ and $\xi_{2} \in[-2,0]$, then we have $C_{p_{2}}(\bar{y})-C_{p_{1}}(\bar{y}) \succeq 0_{\mathbb{R}^{m}}, y_{p_{1}}=u_{p_{1}}$ and $y_{p_{2}}=l_{p_{2}}$. Hence, $\bar{y}=(30,0)$ is the robust vector equilibrium with respect to the worst case. It can be seen that the robust vector equilibrium with respect to the worst case is sensitive to parameter perturbations.

## 3. Min-Max Method for Robust Vector Equilibrium

In this section, a min-max algorithm is proposed to look for a subset of the robust vector equilibrium flows.

### 3.1. Description of the Algorithm

In this subsection, we construct an optimization problem whose solution is equivalent to the the robust vector equilibrium flow. For $(y, \xi) \in \Omega \times U$, we define:

$$
\psi(y, \xi):=\sum_{p_{k}, p_{j} \in P_{\omega}, \omega \in W}\left(y_{p_{k}}-l_{p_{k}}\right)\left(u_{p_{j}}-y_{p_{j}}\right)\left[c_{p_{k}}(y, \xi)-c_{p_{j}}(y, \xi)\right]^{\top} H_{+}\left[c_{p_{k}}(y, \xi)-c_{p_{j}}(y, \xi)\right] .
$$

Proposition 1. Let $\bar{y}$ be a feasible flow. The following statements are equivalent.
(i) $\bar{y}$ is a robust vector equilibrium;
(ii) There exists $(\bar{y}, \bar{\xi})$ such that it is a saddle-point of the problem, denoted as follows:

$$
\begin{array}{r}
\min _{y \in \Omega} \max _{\xi \in U} \psi(y, \xi)  \tag{2}\\
\text { s.t. } y \in \Omega .
\end{array}
$$

and $\psi(\bar{y}, \bar{\xi})$ is equal to zero.
Proof. Firstly, we prove the implication $(i) \Rightarrow(i i)$. Since $\psi(y, \xi) \geq 0$, it suffices to prove $\psi(\bar{y}, \xi)=0$ for all $\xi \in U$, i.e., $0=\psi(\bar{y}, \xi)=\psi(\bar{y}, \bar{\xi}) \leq \psi(y, \xi)$. Hence, for every $\xi \in \mathcal{U}, p_{k} \in P_{\omega}, \omega \in \mathcal{W}$, we consider the following term:

$$
O_{p_{k}}=\sum_{p_{j} \in P_{\omega}}\left(y_{p_{k}}-l_{p_{k}}\right)\left(u_{p_{j}}-y_{p_{j}}\right)\left[c_{p_{k}}(y, \xi)-c_{p_{j}}(y, \xi)\right]^{\top} H_{+}\left[c_{p_{k}}(y, \xi)-c_{p_{j}}(y, \xi)\right] .
$$

If $c_{p_{k}}(\bar{y}, \xi)-c_{p_{j}}(\bar{y}, \xi) \succeq 0_{\mathbb{R}^{m}}$ for some $p_{j} \in P_{\omega}$, then by Definition 2, one has $\bar{y}_{p_{k}}=l_{p_{k}}$ or $\bar{y}_{p_{j}}=u_{p_{j}}$ and so $O_{p_{k}}=0$. If $c_{p_{k}}(\bar{y}, \xi)-c_{p_{j}}(\bar{y}, \xi) \prec 0_{\mathbb{R}^{m}}$, for some $p_{j} \in P_{\omega}, H_{+}\left[c_{p_{k}}(\bar{y}, \xi)-\right.$ $\left.c_{p_{j}}(\bar{y}, \xi)\right]=0_{\mathbb{R}^{m}}$, and hence, $O_{p_{k}}=0$. By the above cases, one has $\psi(\bar{y}, \xi)=0$ for all $\xi \in U$.

Conversely, if (ii) is satisfied, $(\bar{y}, \bar{\zeta})$ is a saddle-point and $O_{p_{k}}=0$. If for every $\xi \in U$, some $p_{k}, p_{j} \in P_{\omega}, \omega \in W$ one has $c_{p_{k}}(\bar{y}, \xi)-c_{p_{j}}(\bar{y}, \xi) \succeq 0_{\mathbb{R}^{m}}$, then $\left[c_{p_{k}}(\bar{y}, \xi)-\right.$ $\left.c_{p_{j}}(\bar{y}, \xi)\right]^{\top} H_{+}\left[c_{p_{k}}(\bar{y}, \xi)-c_{p_{j}}(\bar{y}, \xi)\right]>0$ and so $\bar{y}_{p_{k}}=l_{p_{k}}$ or $\bar{y}_{p_{j}}=u_{p_{j}}$. Consequently, $\bar{y}$ is a robust vector equilibrium.

Now, a min-max algorithm is proposed to solve problem (2). In our algorithm, we select different steps for the two variables $y$ and $\xi$, which is different from one proposed in [22]. In addition, we extend the search directions of the algorithm to make the search faster and more suitable for different needs. Thus, our algorithm is an improvement of that in [22].

Direction set: The set $D$ consist of finite unit vectors which can span $\mathbb{R}^{n}$. Here, in order to reduce the computational cost, we only consider some directions in $D$. For example, when $n=2$, in this paper, let $D=\left\{(1,0),\left(\frac{\sqrt{3}}{2}, \frac{1}{2}\right),(1,1),\left(\frac{1}{2}, \frac{\sqrt{3}}{2}\right),(0,1),\left(-\frac{1}{2}, \frac{\sqrt{3}}{2}\right),(-1,1)\right.$, $\left.\left(-\frac{\sqrt{3}}{2}, \frac{1}{2}\right),(-1,0),\left(-\frac{\sqrt{3}}{2},-\frac{1}{2}\right),(-1,-1),\left(-\frac{1}{2},-\frac{\sqrt{3}}{2}\right),(0,-1),\left(\frac{1}{2},-\frac{\sqrt{3}}{2}\right),(1,-1),\left(\frac{\sqrt{3}}{2},-\frac{1}{2}\right)\right\}$.

Step length: Let initial step $t_{0}=1$ and $d_{k}=\arg \min _{y_{i}} \psi_{k}\left(y_{k}+t_{k} d, \xi_{k}\right), d \in D$. Let $\tilde{y}_{k}=y_{k}+t_{k} d_{k}$. If iteration is successful, i.e., $\left.\psi_{k}\left(\tilde{y}_{k}, \xi_{k}\right)\right)<\psi_{k}\left(y_{k}+t_{k} d, \xi_{k}\right)-c t_{k}^{2}$ for all $d \in D(c>0)$, then the next step length value $t_{k+1}=1$; if the iteration is unsuccessful, then $t_{k+1}=\left\|\tilde{y}-y_{k}\right\| / 2$.

Remark 2. It is worth noting that computations for $y_{t}$ and $\xi_{t}$ in Algorithm 1 are based on Algorithm 2.

```
Algorithm 1: Min-max algorithm (Algorithm 1).
    input : \(\psi\) : objective function; c : forcing function constant \(c>0\);
            \(T\) : maximum number of iterations; \(t_{0}\) : initialize step size;
            \(\left(y_{0}, \xi_{0}\right)\) : initial iteration point; \(S=\varnothing, S E=\varnothing\).
    for \(t=1, \cdots, T\) do
        \(\xi_{t}=\mathbf{A 1}\left(-\psi\left(y_{t-1},.\right), \xi_{t-1}\right)\)
        \(y_{t}=\mathbf{A} \mathbf{1}\left(\psi\left(., \xi_{t}\right), y_{t-1}\right)\)
        return \(\left(y_{T}, \xi_{T}\right)\), store it in \(S\).
    Choose a \(\left(y_{T}, \xi_{T}\right)\) from \(S\), compute \(\psi\left(y_{T}, \xi_{T}\right)\).
    If \(\psi\left(y_{T}, \xi_{T}\right) \leq \epsilon\), store \(y_{T}\) in \(S E\) and return to Step 5 until no element of \(S\) left.
```

```
Algorithm 2: Algorithm \(2\left(\psi(), y_{0}\right)\).
    input \(: \psi\) : objective function; \(c\) : forcing function constant \(c>0\);
            T : maximum number of iterations; \(t_{0}\) : initialize step size value;
            \(y_{0}\) : initial iteration point.
    for \(k=0, \cdots, T-1\) do
        1. Generate direction set
        \(D=\left\{d^{i}\right.\) : any one unit direction of a certain point \(\}\).
        2. Generate the points
\[
y=y_{k}+t_{k} d \subset \Omega, \forall d \in D
\]
3. Choose \(d_{k}=\arg \min _{d} \psi\left(y, \xi_{k}\right)\) and let \(\tilde{y_{k}}=y_{k}+t_{k} d_{k}\).
4. if \(\psi\left(\tilde{y}, \xi_{k}\right)<\psi\left(y_{k}, \xi_{k}\right)-c t_{k}^{2}\) then
(Iteration is successful)
\(y_{k+1}=\tilde{y}, t_{k+1}=1\);
else
(Iteration is unsuccessful)
\(y_{k+1}=y_{k}, t_{k+1}=\left\|\tilde{y}-y_{k}\right\| / 2\).
return \(y_{T}\)
```


### 3.2. Comparison with Other Methods

In this subsection, we will give three numerical examples to show the comparison with that of [22]. In these numerical examples, both Algorithm 1 and the algorithm proposed in [22] start from the same set of initial points. To make a fair comparison, all test problems are run five times to reduce the impact of randomness.

Remark 3. There is a step calculation method in reference [22]-if the iteration is successful: $t_{k+1}=\min \left(t_{\max }, \gamma t_{k}\right), \gamma>1$, where $t_{\max }$ is the largest step size; if iteration is unsuccessful: $t_{k+1}=\frac{1}{\gamma} t_{k}$. Compared with the step calculation method in reference [22], the step calculation method presented in Algorithm 1 has better performance, since Algorithm 1 selects different step size for different variables and extends the search directions. What is more, Algorithm 1 requires neither gradient information nor redundant parameters.

Example 2. Consider a network problem depicted in Figure 1, where $\mathcal{N}=\{1,2\}$, $\mathcal{W}=\{\omega\}=\{(1,2)\}, \mathcal{E}=\left\{\alpha_{1}, \alpha_{2}\right\}, \mathcal{D}=d_{\omega}=30$. There are two criteria: travel time and travel cost. The cost functions of arcs and constraints of paths are given as bellow: $t_{1, \alpha_{1}}(y, \xi)=y_{1}^{2}+2 y_{1} y_{2}+y_{2}-\xi_{1}, t_{2, \alpha_{1}}(y, \xi)=y_{1}+y_{2}^{2} ; t_{1, \alpha_{2}}(y, \xi)=y_{1}^{2}+10 y_{2} y_{2}$, $t_{2, \alpha_{2}}(y, \xi)=7 y_{1}+6 y_{2}^{2}-6 \xi_{2}$.

$$
l_{p_{1}}=0, l_{p_{2}}=0 ; u_{p_{1}}=30, u_{p_{2}}=30
$$



Figure 1. Network topology for Example 2.
Then, we have:

$$
c_{p_{1}}(y, \xi)=\binom{y_{1}^{2}+2 y_{1} y_{2}+y_{2}-\xi_{1}}{y_{1}+y_{2}^{2}}, \quad c_{p_{2}}(y, \xi)=\binom{y_{1}^{2}+10 y_{2} y_{2}}{7 y_{1}+6 y_{2}^{2}-6 \xi_{2}} .
$$

where $\xi_{1} \in[0,1]$ and $\xi_{2} \in[0,1]$. Initial feasible flows and a subset of the robust vector equilibrium flows are obtained in 23.82s. The results are shown in Table 1. However, if we use step calculational method in [22], then it takes 25.68 s and the obtain the same robust vector equilibrium flows with our algorithm.

Table 1. Computational results of Algorithm 1.

| Initial Feasible Flows | Robust Vector Equilibrium Flows |
| :--- | :--- |
| $(0.00,30.00)(3.75,26.25)$ | $(25.00,5.00)(25.75,4.25)$ |
| $(11.25,18.75)(11.25,18.75)$ | $(25.50,4.50)(25.25,4.75)$ |
| $(15.00,15.00)(18.75,11.25)$ | $(25.00,5.00)(25.75,4.25)$ |
| $(22.50,7.50)(26.25,3.75)$ | $(25.50,4.50)(26.25,3.75)$ |
| $(30.00,0.00)$ | $(30.00,0.00)$ |

Example 3. Consider the network problem depicted in Figure 2, where $\mathcal{N}=\{1,2,3,4\}$, $\mathcal{W}=\{(1,4),(2,4)\}$, and there are two $O-D$ pairs, $\omega_{1}, \omega_{2} . \mathcal{E}=\left\{\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}\right\}$, $p_{1}=\left(\alpha_{1} \alpha_{5}\right), p_{2}=\left(\alpha_{2}\right), p_{3}=\left(\alpha_{3} \alpha_{5}\right), p_{4}=\left(\alpha_{4}\right), \mathcal{D}=\left\{d_{\omega_{1}}, d_{\omega_{2}}\right\}=\{55,35\}$. There are two criteria: travel time and travel cost. Constrains of paths are given as follows:

$$
\begin{gathered}
l_{p_{1}}=0, l_{p_{2}}=0, l_{p_{3}}=0, l_{p_{4}}=0 \\
u_{p_{1}}=55, u_{p_{2}}=55, u_{p_{3}}=35, u_{p_{4}}=35
\end{gathered}
$$



Figure 2. Network topology for Example 3.
The cost functions of arcs are defined as follows:

$$
\begin{gathered}
t_{\alpha_{1}}(y, \xi)=\binom{y_{1}^{2}+2 y_{2}+y_{3}^{2}+y_{4}^{2}+1}{2 y_{1}+y_{2}^{2}+2 y_{3}+y_{4}+\frac{3}{2}-6 \xi_{1}}, t_{\alpha_{2}}(y, \xi)=\binom{2 y_{1}+3 y_{2}+5 y_{3}+y_{4}+1+\xi_{2}}{2 y_{1} y_{2}+y_{2}+y_{3}^{2}+y_{4}+\frac{1}{2}}, \\
t_{\alpha_{3}}(y, \xi)=\binom{2 y_{1}+y_{2}+y_{3}^{2}+y_{4}+\frac{3}{2}+3 \xi_{3}}{y_{3}^{2}+5 y_{4}+1}, t_{\alpha_{4}}(y, \xi)=\binom{2 y_{2}+3 y_{3} y_{4}+y_{4}^{2}+1}{y_{1}^{2}+y_{2}+y_{4}+\frac{1}{2}-8 \xi_{4}}, t_{\alpha_{5}}(y, \xi)=\binom{y_{1}+y_{2}}{y_{3}+y_{4}} .
\end{gathered}
$$

Then, we have:

$$
\begin{gathered}
c_{p_{1}}(y, \xi)=\binom{y_{1}^{2}+y_{1}+3 y_{2}+y_{3}^{2}+y_{4}^{2}+1}{2 y_{1}+y_{2}^{2}+3 y_{3}+2 y_{4}+1.5-6 \xi_{1}}, c_{p_{2}}(y, \xi)=\binom{2 y_{1}+3 y_{2}+5 y_{3}+y_{4}+1+\xi_{2}}{2 y_{1} y_{2}+y_{2}+y_{3}^{2}+y_{4}+0.5} \\
\quad c_{p_{3}}(y, \xi)=\binom{3 y_{1}+2 y_{2}+y_{3}^{2}+y_{4}+1.5+3 \xi_{3}}{y_{3}^{2}+y_{3}+6 y_{4}+1}, c_{p_{4}}(y, \xi)=\binom{2 y_{2}+3 y_{3} y_{4}+y_{4}^{2}+1}{y_{1}^{2}+y_{2}+8 y_{4}+0.5-8 \xi_{4}}
\end{gathered}
$$

where $\xi_{i} \in[0,1], i=1,2,3,4$. Initial feasible flows and a subset of the robust vector equilibrium flows are obtained in 40.56 s . The results are shown in Table 2. The time cost of Algorithm 1 is $4 \%$ lower than that of the step calculation method in [22].

Table 2. Computational results of Algorithm 1.

| Initial Feasible Flows | Robust Vector Equilibrium Flows |
| :--- | :--- |
| $(0.00,55.00,0.00,35.00)$ | $(27.00,28.00,27.00,8.00)$ |
| $(5.00,50.00,5.00,30.00)$ | $(27.00,28.00,27.00,8.00)$ |
| $(10.00,45.00,10.00,25.00)$ | $(27.00,28.00,27.00,8.00)$ |
| $(15.00,40.00,15.00,20.00)$ | $(27.00,28.00,27.00,8.00)$ |
| $(20.00,35.00,20.00,15.00)$ | $(27.00,28.00,27.00,8.00)$ |
| $(25.00,30.00,25.00,10.00)$ | $(29.00,26.00,29.00,6.00)$ |
| $(30.00,25.00,30.00,5.00)$ | $(29.00,26.00,29.00,6.00)$ |
| $(35.00,20.00,35.00,0.00)$ |  |

Example 4. Consider the network problem depicted in Figure 3, where $\mathcal{N}=\{1,2,3,4,5,6\}$, $\mathcal{W}=\left\{\omega_{1}, \omega_{2}\right\}=\{(1,5),(2,6)\}, \mathcal{E}=\left\{\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}, \alpha_{6}, \alpha_{7}, \alpha_{8}, \alpha_{9}\right\}, \mathcal{D}=\left\{d_{\omega_{1}}, d_{\omega_{2}}\right\}$, $d_{\omega_{1}}=25, d_{\omega_{2}}=20$, with two criteria: travel time and travel cost. $P_{\omega}=\left\{p_{1}, p_{2}, p_{3}, p_{4}, p_{5}, p_{6}, p_{7}\right\}$, where $P_{\omega_{1}}=\left\{p_{1}, p_{2}, p_{3}, p_{4}\right\}, P_{\omega_{2}}=\left\{p_{5}, p_{6}, p_{7}\right\}, p_{1}=\left(\alpha_{3}\right), p_{2}=\left(\alpha_{2} \alpha_{5} \alpha_{8}\right)$, $p_{3}=\left(\alpha_{1} \alpha_{4} \alpha_{5} \alpha_{8}\right), p_{4}=\left(\alpha_{1} \alpha_{6} \alpha_{8}\right), p_{5}=\left(\alpha_{7}\right), p_{6}=\left(\alpha_{6} \alpha_{9}\right)$, and $p_{7}=\left(\alpha_{4} \alpha_{5} \alpha_{9}\right)$.


Figure 3. Network topology for Example 4.
The constrains of paths and cost functions are given as follows:

$$
\begin{gathered}
l_{p_{1}}=0, l_{p_{2}}=0, l_{p_{3}}=0, l_{p_{4}}=0, l_{p_{5}}=0, l_{p_{6}}=0, l_{p_{7}}=0 \\
u_{p_{1}}=25, u_{p_{2}}=25, u_{p_{3}}=25, u_{p_{4}}=25, u_{p_{5}}=20, u_{p_{6}}=20, u_{p_{7}}=20 \\
t_{1, \alpha_{1}}(y, \xi)=4\left(y_{6}+y_{7}\right)+50-2 \xi_{3}, t_{2, \alpha_{1}}(y, \xi)=\left(y_{6}+y_{7}\right)^{2}+90-2 \xi_{3} ; \\
t_{1, \alpha_{2}}(y, \xi)=2 y_{2}+20-\xi_{2}+4 \xi_{6}, t_{2, \alpha_{2}}(y, \xi)=3 y_{2}^{2}+10 ; \\
t_{1, \alpha_{3}}(y, \xi)=4 y_{1}^{2}+100+\xi_{1}, t_{2, \alpha_{3}}(y, \xi)=2 y_{1}^{2}+110+6 \xi_{1} ; \\
t_{1, \alpha_{4}}(y, \xi)=2\left(y_{4}+y_{7}\right)+10+\xi_{2}, t_{2, \alpha_{4}}(y, \xi)=\left(y_{4}+y_{7}\right)^{2}+30-\xi_{2}+\xi_{3} ; \\
t_{1, \alpha_{5}}(y, \xi)=2\left(y_{3}+y_{4}+y_{7}\right)^{2}+10-\xi_{2}, t_{2, \alpha_{5}}(y, \xi)=\left(y_{3}+y_{4}+y_{7}\right)^{2}+20+\xi_{2} ; \\
t_{1, \alpha_{6}}(y, \xi)=5\left(y_{4}+y_{5}\right)^{2}+430+2 \xi_{3}+\xi_{4}, t_{2, \alpha_{6}}(y, \xi)=2\left(y_{4}+y_{5}\right)+530+2 \xi_{3}-\xi_{4} ; \\
t_{1, \alpha_{7}}(y, \xi)=2 y_{2}^{2}+100+5 \xi_{5}, t_{2, \alpha_{7}}(y, \xi)=3 y_{5}+300 ; \\
t_{1, \alpha_{9}}(y, \xi)=\left(y_{5}+y_{6}\right)^{2}+30-4 \xi_{6}, t_{2, \alpha_{9}}(y, \xi)=2\left(y_{5}+y_{6}\right)+10+2 \xi_{7}-\xi_{3} ;
\end{gathered}
$$

where $z_{\alpha_{i}}(i=1,2, \ldots, 9)$ denotes the flow on arc $\alpha_{i}$. Then we have

$$
\begin{gathered}
c_{p_{1}}(y, \xi)=\binom{4 y_{1}^{2}+100+\xi_{1}}{2 y_{1}^{2}+110+6 \xi_{1}}, c_{p_{2}}(y, \xi)=\binom{2 y_{2}+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{6}+y_{7}\right)^{2}+50-2 \xi_{2}}{3 y_{2}^{2}+\left(y_{3}+y_{4}+y_{7}\right)^{2}+2\left(y_{3}+y_{6}+y_{7}\right)^{2}+40+\xi_{2}} \\
c_{p_{3}}(y, \xi)=\binom{4\left(y_{6}+y_{7}\right)+2\left(y_{4}+y_{7}\right)+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{6}+y_{7}\right)^{2}+90-2 \xi_{3}}{\left(y_{6}+y_{7}\right)^{2}+\left(y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{4}+y_{7}\right)^{2}+2\left(y_{3}+y_{6}+y_{7}\right)^{2}+150-\xi_{3}} \\
c_{p_{4}}(y, \xi)=\binom{4\left(y_{6}+y_{7}\right)+5\left(y_{4}+y_{5}\right)^{2}+\left(y_{3}+y_{6}+y_{7}\right)^{2}+500+\xi_{4}}{\left(y_{6}+y_{7}\right)^{2}+2\left(y_{4}+y_{5}\right)+2\left(y_{3}+y_{6}+y_{7}\right)^{2}+630-\xi_{4}} \\
c_{p_{5}}(y, \xi)=\binom{2 y_{2}^{2}+100+5 \xi_{5}}{3 y_{5}+300}, c_{p_{6}(y, \xi)=\binom{5\left(y_{4}+y_{5}\right)^{2}+\left(y_{5}+y_{6}\right)^{2}+460-4 \xi_{6}}{2\left(y_{4}+y_{5}\right)+2\left(y_{5}+y_{6}\right)+540}}^{c_{p_{7}}(y, \xi)=\binom{2\left(y_{4}+y_{7}\right)+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+\left(y_{5}+y_{6}\right)^{2}+50}{\left(y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{4}+y_{7}\right)^{2}+2\left(y_{5}+y_{6}\right)^{2}+60+2 \xi_{7}}}
\end{gathered}
$$

where $\xi_{i} \in[0,1], i=1,2,3,4,5,6,7$. Initial feasible flows and a subset of the robust vector equilibrium flows are obtained in 538.47 s. The results are shown in Table 3. The algorithm proposed in [22] obtains the same robust vector equilibrium flows, but its time cost is 548.32 s .

Table 3. Computational results of Algorithm 1.

| Initial Feasible Flows | Robust Vector Equilibrium Flows |
| :--- | :--- |
| $(0,0,0,25,0,0,20)(0,0,25,0,0,0,20)$ | $(9,7,9,0,19,0,1)(12,13,0,0,8,0,12)$ |
| $(0,25,0,0,0,0,20)(25,0,0,0,0,0,20)$ | $(15,10,0,0,9,0,11)(15,10,0,0,9,0,11)$ |
| $(0,0,0,25,0,20,0)(0,0,25,0,0,20,0)$ | $(12,8,5,0,15,0,5)(15,10,0,0,9,0,11)$ |
| $(0,25,0,0,0,20,0)(25,0,0,0,0,20,0)$ | $(9,7,9,0,19,0,1)(12,13,0,0,8,0,12)$ |
| $(0,0,0,25,20,0,0)(0,0,25,0,20,0,0)$ | $(11,7,7,0,18,0,2)$ |
| $(0,25,0,0,20,0,0)(25,0,0,0,20,0,0)$ |  |

## 4. Smoothing Method for the Robust Vector Equilibrium with the Worst Case

It is worth noting that the algorithm in [21] needs to solve a non-smoothing optimization problem. This results in its computational inefficiency. This prompts us to continuously investigate algorithm for solving robust equilibrium flows. In this section, we propose a smoothing method to calculate a subset of the robust vector equilibrium with respect to the worst case. The algorithm is denoted Algorithm 3. To generalize a subset of the robust vector equilibrium flows with respect to the worst case, we use a two-step strategy. The first step is to construct an equivalent optimization problem with the help of a variant version of ReLU activation function for finding the robust weak vector equilibrium flows with respect to the worst case. The second step is to judge whether or not the robust weak vector equilibrium flows with respect to the worst case are equal to the robust vector equilibrium flows with respect to the worst case by an equivalent optimization problem using the vector version of heaviside step function.

```
Algorithm 3: Robust vector equilibrium algorithm (denoted Algorithm 3).
    1 Choose a positive integer q and a tolerance level \(\epsilon \geq 0\).
    2 Enter \(l=\left(l_{p_{k}}\right)_{p_{k} \in P}\) and \(u=\left(u_{p_{k}}\right)_{P_{k} \in P}\). Set \(\delta_{j}=d_{\omega_{j}} /\left(q\left|P_{\omega_{j}}\right|\right), j=1, \cdots, \tilde{l}\).
    \({ }_{3}\) Choose \(\left(k_{1}, \cdots, k_{n}\right)^{\top} \in \mathbb{N}^{n}\) satisfying
\[
\sum_{i \in I_{j}} k_{i}=q\left|P_{\omega_{j}}\right|, \text { and } l_{p_{i}} \leq k_{i} \delta_{j} \leq u_{p_{i}}, i \in I_{j}, j=1, \cdots, \tilde{l}
\]
4 Store \(y=\left(y_{p_{1}}, \cdots, y_{p_{n}}\right)^{\top}\) in \(S^{0}\) where
\[
y_{p_{i}}=k_{i} \delta_{j}, i \in I_{j}, j=1, \cdots, \tilde{l}
\]
and return to Step 3 for other vectors ( \(k_{1}, \cdots, k_{n}\) ) unless no one left.
5 Choose a feasible flow \(y^{0}\) from \(S^{0}\) to start. Set \(k=0, S^{0}=S^{0} \backslash\left\{y^{0}\right\}\) and \(W E=\varnothing\).
6 For every \(i, j \in\{1, \cdots, n\}\), solve
\[
\begin{aligned}
& \operatorname{minimize} \\
& \text { subject to } \\
& \\
& \qquad \mid y \in \Omega \\
& \\
& \left|y_{p_{i}}-y_{p_{i}}^{0}\right| \leq \delta_{\omega(i)}, i=1, \cdots, n
\end{aligned}
\]
```

If $\phi(y) \leqq \epsilon$, store $y$ in $W E$ and return to Step 5 until no element of $S^{0}$ left.
7 Choose a feasible flow $y \in W E, W E=W E \backslash\{y\}$.
8 Compute

$$
\varphi(y)=\sum_{\omega \in W} \sum_{p_{k}, p_{j} \in P_{\omega}}\left(y_{p_{k}}-l_{p_{k}}\right)\left(u_{p_{j}}-y_{p_{j}}\right)\left(C_{p_{k}}(y)-C_{p_{j}}(y)\right)^{\top} H_{+}\left[C_{p_{k}}(y)-C_{p_{j}}(y)\right]
$$

9 If $\varphi(y) \leqq \epsilon$, store $y$ in $E$ and return to Step 7 until no element of $W E$ left.

Define a function $r: \mathbb{R} \rightarrow \mathbb{R}$ and give its vector version function $\mathcal{R}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ as follows:

$$
\begin{aligned}
& r(a)=(\max \{0, a\})^{2} . \\
& \mathcal{R}(x)=\left(\prod_{i=1}^{n} r\left(x_{i}\right)\right) e
\end{aligned}
$$

In addition, the heaviside step function $h_{+}: \mathbb{R} \rightarrow \mathbb{R}$ and its vector version function $H_{+}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ are also given below:

$$
\begin{gathered}
h_{+}(a)= \begin{cases}1, & \text { if } a \geq 0, \\
0, & \text { otherwise } .\end{cases} \\
H_{+}(x)=\left(\prod_{i=1}^{n} h_{+}\left(x_{i}\right)\right) e, \forall x \in \mathbb{R}^{m} .
\end{gathered}
$$

### 4.1. Description of the Algorithm

In this subsection, we construct an optimization problem whose solution is equivalent to a robust weak vector equilibrium flow with respect to the worst case. For $y \in \Omega$, we define:

$$
\phi(y):=\sum_{\omega \in W} \sum_{p_{k}, p_{j} \in P_{\omega}}\left(y_{p_{k}}-l_{p_{k}}\right)\left(u_{p_{j}}-y_{p_{j}}\right)\left(C_{p_{k}}(y)-C_{p_{j}}(y)\right)^{\top} \mathcal{R}\left[C_{p_{k}}(y)-C_{p_{j}}(y)\right] .
$$

Proposition 2. Let $\bar{y}$ be a feasible flow. The following statements are equivalent.
(i) $\bar{y}$ is a robust weak vector equilibrium with respect to the worst case;
(ii) $\bar{y}$ is an optimal solution of the problem, denoted:

$$
\begin{array}{r}
\min \phi(\bar{y}) \\
\text { s.t. } y \in \Omega . \tag{3}
\end{array}
$$

and the optimal value $\phi(\bar{y})$ is equal to zero.
Proof. We first prove the implication $(i) \Rightarrow(i i)$. It is not hard to see $\phi(y) \geq 0$ for every $y \in \Omega$. Thus, if $\bar{y}$ is a robust weak vector equilibrium with respect to the worst case, in order to deduce (ii), it suffices to prove $\phi(\bar{y})=0$. In addition, for every $p_{k} \in p_{\omega}, \omega \in W$, consider the term:

$$
Q_{p}=\sum_{\omega \in \mathcal{W}} \sum_{p_{k}, p_{j} \in P_{\omega}}\left(y_{p_{k}}-l_{p_{k}}\right)\left(u_{p_{j}}-y_{p_{j}}\right)\left(C_{p_{k}}(y)-C_{p_{j}}(y)\right)^{\top} \mathcal{R}\left[C_{p_{k}}(y)-C_{p_{j}}(y)\right] .
$$

If $C_{p_{k}}(\bar{y})-C_{p_{j}}(\bar{y}) \succ 0_{\mathbb{R}^{m}}$, for some $p_{j} \in P_{\omega}$, then by Definition 4, either $\bar{y}_{p_{k}}=l_{p_{k}}$ or $\bar{y}_{p_{j}}=u_{p_{j}}$; if $C_{p_{k}}(\bar{y})-C_{p_{j}}(\bar{y})=0_{\mathbb{R}^{m}}$, for some $p_{j} \in P_{\omega}$, we also get $Q_{p}=0$; if $C_{p_{k}}(\bar{y})-C_{p_{j}}(\bar{y}) \prec 0_{\mathbb{R}^{m}}$, then $\mathcal{R}\left[C_{p_{k}}(y)-C_{p_{j}}(y)\right]=0_{\mathbb{R}^{m}}$, and thus, $Q_{p}=0$. As a result, one has $\phi(\bar{y})=0$.

Conversely, assume that $\bar{y}$ is an optimal solution of Problem (3) and $\phi(\bar{y})=0$. Then, we have $Q_{p}=0$ for all $p \in P$. If there exists some $p_{k}, p_{j} \in p_{\omega}, \omega \in \mathcal{W}$ such that $C_{p_{k}}(\bar{y})-C_{p_{j}}(\bar{y}) \succ 0_{\mathbb{R}^{m}}$, then $\left(C_{p_{k}}(y)-C_{p_{j}}(y)\right)^{\prime} \mathcal{R}\left[C_{p_{k}}(y)-C_{p_{j}}(y) \succ 0_{\mathbb{R}^{m}}\right.$, and thus, either $\bar{y}_{p_{k}}=l_{p_{k}}$ or $\bar{y}_{p_{j}}=u_{p_{j}}$ by $Q_{p}=0$. Consequently, we deduce that $\bar{y}$ is a robust weak vector equilibrium with respect to the worst case.

For $y \in \Omega$, we define

$$
\varphi(y)=\sum_{\omega \in \mathcal{W}} \sum_{p_{k}, p_{j} \in P_{\omega}}\left(y_{p_{k}}-l_{p_{k}}\right)\left(u_{p_{j}}-y_{p_{j}}\right)\left(C_{p_{k}}(y)-C_{p_{j}}(y)\right)^{\top} H_{+}\left[C_{p_{k}}(y)-C_{p_{j}}(y)\right] .
$$

Then, by using a similar method of proof, we may establish the following result for the robust vector equilibrium with respect to the worst case.

Proposition 3. Let $\bar{y}$ be a feasible flow. The following statements are equivalent.
(i) $\bar{y}$ is a robust vector equilibrium with respect to the worst case.
(ii) $\bar{y}$ is an optimal solution of the problem, denoted as follows:

$$
\begin{gather*}
\min \varphi(y)  \tag{4}\\
\text { s.t. } y \in \Omega .
\end{gather*}
$$

and the optimal value $\varphi(\bar{y})$ is equal to zero.
Algorithm 3 is mainly based on ideas of Propositions 2 and 3. Steps 1-4 create a subset of feasible flows with the initial conditions, denoted as $S^{0}$, with which Steps 4-6 will start. Steps 5-6 are aimed at solving Problem (3) given in Proposition 2 by using first-order optimization methods, and then a subset of the robust weak vector equilibrium flows with respect to the worst case is gained. Steps 7-9 focus on solving Problem (4) given in Proposition 3, and then a subset of the robust vector equilibrium flows with respect to the worst case is generated.

Assume that $\mathcal{W}$ consists of $\tilde{l}$ elements $\omega_{1}, \ldots, \omega_{\tilde{l}}$ in the network and for each pair $\omega_{i}$. Let $I_{j}=\left\{i \in\{1, \ldots, n\}: p_{i} \in P_{w_{j}}\right\}$. Denote $W E$ by the subset of the robust weak vector equilibrium flows with respect to the worst case and $E$ by the subset of the robust vector equilibrium flows with respect to worst case.

### 4.2. Comparison with Other Methods

In this subsection, we will give two numerical examples to show the comparison with that of [21]. In these numerical examples, both Algorithm 3 and the algorithm proposed in [21] start from the same set of initial points. To make a fair comparison, all test problems are run five times to reduce the impact of randomness.

Example 5. Consider the network problem depicted in Figure 4, where $\mathcal{N}=\{1,2,3,4,5\}$, $\mathcal{W}=\left\{\omega_{1}, \omega_{2}\right\}=\{(1,4),(1,5)\}, \mathcal{E}=\left\{\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}, \alpha_{6}, \alpha_{7}, \alpha_{8}, \alpha_{9}\right\}, \mathcal{D}=\left\{d_{\omega_{1}}, d_{\omega_{2}}\right\}$, $d_{\omega_{1}}=25, d_{\omega_{2}}=20$, with two criteria: travel time and travel cost, $P_{\omega}=\left\{p_{1}, p_{2}, p_{3}, p_{4}, p_{5}, p_{6}, p_{7}\right\}$, where $P_{\omega_{1}}=\left\{p_{1}, p_{2}, p_{3}, p_{4}\right\}, P_{\omega_{2}}=\left\{p_{5}, p_{6}, p_{7}\right\}$.


Figure 4. Network topology for Example 5.
Assume that:

$$
\begin{gathered}
l_{p_{1}}=l_{p_{2}}=l_{p_{3}}=l_{p_{4}}=l_{p_{5}}=l_{p_{6}}=l_{p_{7}}=l_{p_{8}}=0 \\
u_{p_{1}}=100, u_{p_{2}}=100, u_{p_{3}}=100, u_{p_{4}}=200, u_{p_{5}}=100, u_{p_{6}}=120, u_{p_{7}}=150, u_{p_{8}}=100 ;
\end{gathered}
$$

$$
\begin{gathered}
t_{1, \alpha_{1}}(y, \xi)=y_{1}^{2}+y_{2}^{2}+y_{3}^{3}+\xi_{1}, t_{2, \alpha_{1}}(y, \xi)=2 y_{1}+5 y_{2}+3 y_{3}+y_{4}-2 \xi_{1} ; \\
t_{1, \alpha_{2}}(y, \xi)=8 y_{1} y_{2}+y_{2}^{2}+y_{7}+y_{8}+4 \xi_{2}, t_{2, \alpha_{2}}(y, \xi)=y_{2}+10 y_{3}+2 y_{7}+y_{8}+2 \xi_{2} ; \\
t_{1, \alpha_{3}}(y, \xi)=y_{1}+y_{2}^{2}+y_{3}^{3}+y_{5}+y_{6}-3 \xi_{3}, t_{2, \alpha_{3}}(y, \xi)=10 y_{3}^{3}+2 y_{5}+\xi_{3} ; \\
t_{1, \alpha_{4}}(y, \xi)=y_{1}+y_{2}+y_{4}^{3}+y_{5}^{2}+y_{8}^{3}+\xi_{4}, t_{2, \alpha_{4}}(y, \xi)=y_{1}+2 y_{4}+y_{6} y_{5}+15 y_{8} ; \\
t_{1, \alpha_{5}}(y, \xi)=y_{1}+y_{3}+y_{4}^{3}+y_{5}^{2}+y_{6}^{2}, t_{2, \alpha_{5}}(y, \xi)=y_{1}+5 y_{3}+5 y_{5}+3 y_{6}+12 y_{7}+4 \xi_{5} ; \\
t_{1, \alpha_{6}}(y, \xi)=y_{3}+y_{4}+y_{5}+y_{6}^{3}-3 \xi_{6}, t_{2, \alpha_{6}}(y, \xi)=3 y_{3}+10 y_{5}+y_{6}+2 y_{8}-2 \xi_{6} ; \\
t_{1, \alpha_{7}}(y, \xi)=y_{2}+8 y_{4}^{2}+y_{5}+y_{7}^{3}+2 \xi_{7}, t_{2, \alpha_{7}}(y, \xi)=y_{1}+y_{2}+5 y_{4}+3 y_{7} ; \\
t_{1, \alpha_{8}}(y, \xi)=y_{1}+y_{3}+8 y_{6} y_{7}+y_{8}^{2}+\xi_{8}, t_{2, \alpha_{8}}(y, \xi)=y_{1}+y_{3}+10 y_{5}^{3}+y_{8}-\xi_{8},
\end{gathered}
$$

where $\xi_{i} \in[0,1], i=1,2,3,4,5,6,7,8$. Then, we have

$$
\begin{gathered}
C_{p_{1}}(y)=\binom{y_{1}^{2}+y_{2}^{2}+y_{3}^{3}+1}{2 y_{1}+5 y_{2}+3 y_{3}+y_{4}}, C_{p_{2}}(y)=\binom{8 y_{1} y_{2}+y_{2}^{2}+y_{7}+y_{8}+4}{y_{2}+10 y_{3}+2 y_{7}+y_{8}+2} \\
C_{p_{3}}(y)=\binom{y_{1}+y_{2}^{2}+y_{3}^{3}+y_{5}+y_{6}}{10 y_{3}^{3}+2 y_{5}+1}, C_{p_{4}}(y)=\binom{y_{1}+y_{2}+y_{4}^{3}+y_{5}^{2}+y_{8}^{3}+1}{y_{1}+2 y_{4}+y_{6} y_{5}+15 y_{8}} \\
C_{p_{5}}(y)=\binom{y_{1}+y_{3}+y_{4}^{3}+y_{5}^{2}+y_{6}^{2}}{y_{1}+5 y_{3}+5 y_{5}+3 y_{6}+12 y_{7}+4}, C_{p_{6}}(y)=\binom{y_{3}+y_{4}+y_{5}+y_{6}^{3}}{3 y_{3}+10 y_{5}+y_{6}+2 y_{8}} \\
C_{p_{7}}(y)=\binom{y_{2}+8 y_{4}^{2}+y_{5}+y_{7}^{3}+2}{y_{1}+y_{2}+5 y_{4}+3 y_{7}}, C_{p_{8}}(y)=\binom{y_{1}+y_{3}+8 y_{6} y_{7}+y_{8}^{2}+1}{y_{1}+y_{3}+10 y_{5}^{3}+y_{8}}
\end{gathered}
$$

Choosing $q=2$, we have 32 feasible flows and 2 robust (weak) vector equilibrium flows with respect to the worst case, which are obtained in 0.18 s . Robust (weak) vector equilibrium flows with respect to the worst case are shown in Table 4. However, using the algorithm proposed in [21], it will take $13.826 s$ to obtain five robust vector equilibrium flows with respect to the worst case, which are shown in Table 5.

Table 4. Computational results of Algorithm 3.

| Robust Weak Vector Equilibrium Flows <br> (Worst Case) | Robust Vector Equilibrium Flows (Worst <br> Case) |
| :--- | :--- |
| $(100,100,100,145.125,0,120,134.875,100)$ | $(100,100,100,145.125,0,120,134.875,100)$ |
| $(100,100,100,150,0,120,130,100)$ | $(100,100,100,150,0,120,130,100)$ |

Table 5. Computational results of algorithm in [21].

| Robust Vector Equilibrium Flows (Worst Case) |
| :--- |
| $(100,100,100,30,100,120,150,100)$ |
| $(100,100,100,50,100,120,130,100)$ |
| $(100,100,100,100,100,120,80,100)$ |
| $(100,100,100,150,100,120,30,100)$ |
| $(100,100,100,200,100,100,100,0)$ |

Example 6. Consider the network problem depicted in Figure 5, where $\mathcal{N}=\{1,2,3,4,5\}$, $\mathcal{W}=\left\{\omega_{1}, \omega_{2}\right\}=\{(1,4),(1,5)\}, \mathcal{E}=\left\{\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}, \alpha_{6}, \alpha_{7}, \alpha_{8}, \alpha_{9}\right\}, \mathcal{D}=\left\{d_{\omega_{1}}, d_{\omega_{2}}\right\}$ with $d_{\omega_{1}}=25, d_{\omega_{2}}=20$, with two criteria: travel time and travel cost, $P_{\omega}=\left\{p_{1}, p_{2}, p_{3}, p_{4}, p_{5}, p_{6}, p_{7}\right\}$,
where $P_{\omega_{1}}=\left\{p_{1}, p_{2}, p_{3}, p_{4}\right\}, P_{\omega_{2}}=\left\{p_{5}, p_{6}, p_{7}\right\} . p_{1}=\left(\alpha_{4}\right), p_{2}=\left(\alpha_{2} \alpha_{8}\right), p_{3}=\left(\alpha_{2} \alpha_{6} \alpha_{3}\right)$, $p_{4}=\left(\alpha_{1} \alpha_{3}\right), p_{5}=\left(\alpha_{5}\right), p_{6}=\left(\alpha_{2} \alpha_{9}\right), p_{7}=\left(\alpha_{1} \alpha_{7}\right)$. Let:

$$
\begin{gathered}
l_{1}=0, l_{2}=0, l_{3}=0, l_{4}=0, l_{5}=0, l_{6}=0, l_{7}=0 \\
u_{1}=15, u_{2}=20, u_{3}=15, u_{4}=10, u_{5}=15, u_{6}=10, u_{7}=15 .
\end{gathered}
$$



Figure 5. Network topology for Example 6.
Now, we give the cost function of arcs as follows:

$$
\begin{gathered}
t_{\alpha_{1}}(y, \xi)=\binom{4\left(y_{6}+y_{7}\right)+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+300-2 \xi_{1}}{\left(y_{6}+y_{7}\right)^{2}+\left(y_{3}+y_{4}+y_{7}\right)^{2}+330-\xi_{1}}, t_{\alpha_{2}}(y, \xi)=\binom{2\left(y_{3}+y_{4}+y_{7}\right)^{2}+50+\xi_{2}}{\left(y_{3}+y_{4}+y_{7}\right)^{2}} \\
t_{\alpha_{3}}(y, \xi)=\binom{2\left(y_{3}+y_{7}\right)+2\left(y_{4}+y_{7}\right)^{2}+300+2 \xi_{1}+\xi_{6}}{\left(y_{3}+y_{7}\right)^{2}+2\left(y_{4}+y_{7}\right)+100+\xi_{1}-\xi_{3}}, t_{\alpha_{4}}\left(y_{,} \xi\right)=\binom{4 y_{1}^{2}+100+\xi_{4}}{2 y_{1}^{2}+110+6 \xi_{4}} \\
t_{\alpha_{5}}(y, \xi)=\binom{2 y_{2}^{2}+100+5 \xi_{5}}{3 y_{2}+y_{5}+300}, t_{\alpha_{6}}(y, \xi)=\binom{\left(y_{3}+y_{6}+y_{7}\right)^{2}-260-\xi_{6}-\xi_{2}}{2\left(y_{3}+y_{6}+y_{7}\right)^{2}+50-\xi_{1}} \\
t_{\alpha_{7}}(y, \xi)=\binom{2\left(y_{3}+y_{7}\right)+2\left(y_{4}+y_{7}\right)^{2}-150+3 \xi_{1}-\xi_{7}}{\left(y_{3}+y_{7}\right)^{2}+\left(y_{4}+y_{7}\right)^{2}-310+\xi_{1}}, t_{\alpha_{8}}(y, \xi)=\binom{2 y_{2}+\left(y_{3}+y_{6}+y_{7}\right)^{2}-\xi_{2}}{3 y_{2}^{2}+2\left(y_{3}+y_{6}+y_{7}\right)^{2}+40-\xi_{8}} \\
t_{\alpha_{9}}(y, \xi)=\binom{5\left(y_{5}+y_{6}\right)^{2}+410+\xi_{2}}{2\left(y_{5}+y_{6}\right)^{2}+540-\xi_{9}}
\end{gathered}
$$

where $\xi_{i} \in[0,1], i=1,2,3,4,5,6,7,8,9$. Then, we have:

$$
\begin{gathered}
C_{p_{1}}(y)=\binom{4 y_{1}^{2}+101}{2 y_{1}^{2}+116}, C_{p_{2}}(y)=\binom{2 y_{2}+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{6}+y_{7}\right)^{2}+50}{3 y_{2}^{2}+\left(y_{3}+y_{4}+y_{7}\right)^{2}+2\left(y_{3}+y_{6}+y_{7}\right)^{2}+41}, \\
C_{p_{3}}(y)=\binom{4\left(y_{3}+y_{7}\right)+2\left(y_{4}+y_{7}\right)^{2}+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{6}+y_{7}\right)^{2}+90}{\left(y_{3}+y_{7}\right)^{2}+2\left(y_{4}+y_{7}\right)+\left(y_{3}+y_{4}+y_{7}\right)^{2}+2\left(y_{3}+y_{6}+y_{7}\right)^{2}+150}, \\
C_{p_{4}}(y)=\binom{4\left(y_{3}+y_{7}\right)+2\left(y_{4}+y_{7}\right)^{2}+4\left(y_{6}+y_{7}\right)+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+601}{\left(y_{3}+y_{7}\right)^{2}+2\left(y_{4}+y_{7}\right)+\left(y_{6}+y_{7}\right)^{2}+\left(y_{3}+y_{4}+y_{7}\right)^{2}+430},
\end{gathered}
$$

$$
\begin{gathered}
C_{p_{5}}(y)=\binom{2 y_{2}^{2}+105}{3 y_{2}+y_{5}+300}, C_{p_{6}}(y)=\binom{5\left(y_{5}+y_{6}\right)^{2}+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+460}{2\left(y_{5}+y_{6}\right)+\left(y_{3}+y_{4}+y_{7}\right)^{2}+540}, \\
C_{p_{7}}(y)=\binom{4\left(y_{6}+y_{7}\right)+2\left(y_{4}+y_{7}\right)^{2}+2\left(y_{3}+y_{4}+y_{7}\right)^{2}+2\left(y_{3}+y_{7}\right)+152}{\left(y_{6}+y_{7}\right)^{2}+\left(y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{4}+y_{7}\right)^{2}+\left(y_{3}+y_{7}\right)^{2}+20} .
\end{gathered}
$$

Choosing $q=1,80$ feasible flows are created and a subset of robust vector equilibrium flows with respect to the worst case (displayed in Tables 6) are obtained. This takes about 0.62 s . However, using the algorithm presented in [21], it spends 579 s to obtain a subset of robust vector equilibrium flows with respect to the worst case, which are shown in Table 7.

Table 6. Computational results of Algorithm 3.

| Robust <br> (Worst Case) | Vector Equilibrium Flows |
| :--- | :--- |
| $(11.88,11.29,1.83,0,13.23,0,6.77)$ | Robust Vector Equilibrium Flows (Worst <br> Case) |
| $(10.56,12.5,1.94,0,13.33,0,6.67)$ | $(11.88,11.29,1.83,0,13.23,0,6.77)$ |
| $(10.04,14.92,0.04,0,11.89,0,8.11)$ | $(10.56,12.5,1.94,0,13.33,0,6.67)$ |
| $(11.03,12.4,1.57,0,13.33,0,6.67)$ | $(10.04,14.92,0.04,0,11.89,0,8.11)$ |
| $(11.42,13.58,0,0,11.93,0,8.07)$ | $(11.03,12.4,1.57,0,13.33,0,6.67)$ |
| $(10.45,12.5,2.05,0,13.33,0,6.67)$ | $(11.42,13.58,0,0,11.93,0,8.07)$ |
| $(11.12,10,3.88,0,15,0,5)$ | $(10.45,12.5,2.05,0,13.33,0,6.67)$ |

Table 7. Computational results of algorithm in [21].

```
Robust Vector Equilibrium Flows (Worst Case)
(10.94, 12.5, 1.56, 0, 12.29, 0, 7.71)
(10.94,13.28, 0.78, 0, 12.29, 0,7.71)
(12.66, 10.78, 1.56, 0, 12.29, 0,7.71)
(10.94, 12.5, 1.56, 0, 12.5, 0, 7.5)
```


## 5. Conclusions

In this paper, we mainly consider a robust multi-criteria traffic network equilibrium problem with path capacity constraints. Firstly, the robust vector equilibrium principle and the robust vector equilibrium principle with respect to the worst case are given. We pay attention to constructing an equivalent min-max optimization problem for the robust vector equilibrium, in which the solution is equivalent to a robust vector equilibrium. Then, a direct search algorithm is proposed for solving the corresponding min-max optimization problem. The step size in the algorithm requires neither gradient information nor redundant parameters. What is more, we select different step sizes for different variables and extend the search directions. The results of three numerical experiments show that it takes less time than the method in [22] to find the robust vector equilibrium flows.

To generate a subset of the robust vector equilibrium with respect to the worst case, we employ a two-step strategy. The first step is to construct a smoothing optimization problem based on a variant version of the ReLU activation function to compute the robust weak vector equilibrium flows with respect to the worst case. The second step is to judge whether or not the robust weak vector equilibrium flows with respect to the worst case are equal to the robust vector equilibrium flows with respect to the worst case. Compared with the algorithm in [21], the results of two numerical experiments show that our algorithm can greatly reduce the computational cost.

Recently, robust vector optimization based on set orders is widely used in the uncertain optimization environment $[23,24]$. It is noteworthy that the robust vector equilibrium principles considered in this paper are all based on vector order. In addition, our method can only be applied to small-scale traffic networks. Therefore, an interesting topic for future research is to investigate large-scale, multi-criteria traffic networks based on set orders.

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Article

# A Nonconstant Gradient Constrained Problem for Nonlinear Monotone Operators 

Sofia Giuffrè (1)

Department of Information Engineering, Infrastructure and Sustainable Energy (DIIES) Mediterranea University of Reggio Calabria, Loc. Feo di Vito, 89122 Reggio Calabria, Italy; sofia.giuffre@unirc.it


#### Abstract

The purpose of the research is the study of a nonconstant gradient constrained problem for nonlinear monotone operators. In particular, we study a stationary variational inequality, defined by a strongly monotone operator, in a convex set of gradient-type constraints. We investigate the relationship between the nonconstant gradient constrained problem and a suitable double obstacle problem, where the obstacles are the viscosity solutions to a Hamilton-Jacobi equation, and we show the equivalence between the two variational problems. To obtain the equivalence, we prove that a suitable constraint qualification condition, Assumption S, is fulfilled at the solution of the double obstacle problem. It allows us to apply a strong duality theory, holding under Assumption S. Then, we also provide the proof of existence of Lagrange multipliers. The elements in question can be not only functions in $L^{2}$, but also measures.


Keywords: variational inequalities; non-constant gradient constraints; obstacle problem; nonlinear monotone operators; Lagrange multipliers

MSC: 35J87; 65K10; 49N15

## 1. Introduction

A very interesting problem, which has attracted much interest for many decades because of its simple formulation in terms of differential equations, is the elastic-plastic torsion problem, namely, the problem of minimizing the functional

$$
\begin{equation*}
\frac{1}{2} \int_{\Omega}\left(|D v|^{2}-h v\right) d x \tag{1}
\end{equation*}
$$

on the class of functions $\left\{v \in H_{0}^{1}(\Omega):|D v| \leq 1\right\}$.
The elastic-plastic torsion problem arises when a long elastic bar with cross section $\Omega$ is twisted by an angle. In particular, the formulation due to R. von Mises (see [1]) of the elastic-plastic torsion problem of a cylindrical bar is the following one:
"Find a function $u(x)$, which vanishes on $\partial \Omega$ and is continuous, together with its first derivatives on $\Omega$; on $\Omega$ the gradient of $u$, $D u$, must have an absolute value less than or equal to a given positive constant $t$; whenever, in $\Omega,|D u|<t$, the function $u$ must satisfy the differential equation $\Delta u=-2 v \alpha$, where the positive constants $v$ and $\alpha$ denote the shearing modulus and the angle of twist per unit length respectively".

The plastic region, $P$, refers to the range of deformation in which the material exhibits significant plastic or irreversible behavior. It is the region beyond the elastic limit where the material undergoes permanent changes in shape, and the deformation is not recoverable. When a material is loaded within its elastic limit, it deforms elastically, meaning that it can return to its original shape once the load is removed. However, beyond the elastic limit, the material enters the plastic region, and plastic deformation occurs. In particular, the set

$$
E=\{x \in \Omega:|D u(x)|<t\}
$$

is the set of points where the cross section still remains elastic, namely the elastic set, and the set

$$
P=\{x \in \Omega:|D u(x)|=t\},
$$

is the set of points where the material has become plastic due to the torsion, namely the plastic set.

The ridge $R$ of $\Omega$ is, by definition, the set of points in $\Omega$ where $\operatorname{dist}(x, \partial \Omega)$ is not $C^{1,1}$, whereas the part of $\partial E$, which is contained in $\Omega$, is called the free boundary (see [2]).

For the derivation of the variational inequality from the physical problem see [3].
Ting [4] investigated problem (1) for $n=2$, whereas the existence of a Lagrange multiplier formulation for (1) (and hence of a corresponding system of partial differential equations) was proved for constant $h$ in [5] by Brézis.

Glowinski et al. [6] studied the numerical aspects; for results on the elastic and plastic sets $E$ and $P$ and on the free boundary we refer to Caffarelli and Friedman [2]. In [7] Brezis and Sibony proved that the elastic-plastic torsion problem is equivalent to an obstacle-type problem, in which the distance function represents the obstacle. Moreover, they proposed two numerical methods for the obstacle problem.

In [8] Chiadò Piat and Percivale proved the existence of measure-type Lagrange multipliers under more general assumptions on the operator and on $h$.

Daniele et al. [9] obtained similar results, solving a problem unsolved for a long time by using a new infinite dimensional duality theory. They show, for a class of problems including Problem (1), the existence of an $L^{\infty}$ Lagrange multiplier, if the problem admits solution and a constraint qualification condition is fulfilled at this solution (see Section 3). The Lagrange multiplier is the solution to a dual problem (see also [10-12] for other results related to linear and nonlinear monotone operators).

Many other studies in the past years are related to the problem when the gradient constraints are no longer constant, since it models many interesting physical and biological phenomena (see [13] for an overview of constrained and unconstrained free boundary problems).

Studying variational problems with gradient constraints involves techniques from the calculus of variations and constrained optimization. Some common methods include: Lagrange multipliers, penalty methods, augmented Lagrangian methods, and projection methods. As is well known, Lagrange multipliers introduce additional unknowns and allow the constraints to be incorporated into the objective function through a modified Lagrangian. The resulting problem can then be solved using variational methods or numerical optimization techniques.

Relevant issues related to the problem with gradient constraints are existence and regularity of the solution, existence of Lagrange multipliers, connection with double obstacle problem, and numerical aspects.

Regarding the existence and regularity of the solution, we refer to L. Evans in [14], who studied general linear elliptic equations with a non-constant gradient constraint $g(x) \in$ $C^{2}(\bar{\Omega})$, and proved that there exists a unique solution in the space $W_{l o c}^{2, p}(\Omega) \cap W_{0}^{1, \infty}(\Omega)$, with $1<p<\infty$ (see also [15-18] for other regularity results).

The conditions required for the existence of a Lagrange multiplier are typically related to the regularity of the problem, such as the smoothness of the objective function and constraints.

One of the important conditions for the existence of a Lagrange multiplier is the constraint qualification. There are different types of constraint qualifications, such as the linear independence constraint qualification (LICQ), the Mangasarian-Fromovitz constraint qualification (MFCQ), and Slater's condition (see Theorem 4).

If the qualification condition is satisfied, then according to the Lagrange multiplier theorem, there exists a Lagrange multiplier associated with the optimal solution. The Lagrange multiplier, generally, provides information about the sensitivity of the objective function to changes in the constraints.

However, let us stress that the existence of a Lagrange multiplier does not guarantee a unique solution to the optimization problem. It only indicates the existence of a necessary condition for an optimum.

A very interesting property of variational problems with gradient constraints is the relationship with double obstacle problems. Let us remark that the equivalence is not always true, as observed in [14] (see also [19,20]). Equivalence results between the two problems associated to the Laplacian or to a linear operator are contained in [19-21] (see also [22]).

The equivalence also holds for the problem associated to a nonlinear strongly monotone operator $a(D u)$ with nonconstant gradient constraint of type $G(D u) \leq M$, where $G$ is a strictly convex function (see [23]).

Let us note that monotone operators play a fundamental role in various branches of mathematics, including optimization theory, to analyze several mathematical problems involving nonlinear operators. Monotone operators are extensively used in the study of variational inequalities too (see [24]).

The paper adds to the literature on nonconstant gradient constrained problem further results related to the relationship with double obstacle problem and the existence of Lagrange multipliers. Here we investigate the problem associated to a nonlinear strongly monotone operator as in [23], but we consider the nonconstant gradient constraint of type $|D u| \leq g(x)$, with $g(x) \in C^{2}(\bar{\Omega}), g(x)>0$. We also prove the existence of $L^{2}$ Lagrange multipliers and, under less restrective assumptions, an existence result of measure-type Lagrange multipliers. Let us note that the existence of Lagrange multipliers as measures is not proved for gradient contraints of type $G(D u) \leq M$.

In particular, the problem under consideration is

$$
\begin{align*}
& \text { Find } u \in K_{g}=\left\{v \in H_{0}^{1,2}(\Omega):|D v|^{2}=\sum_{i=1}^{n}\left(D_{i} v\right)^{2} \leq g(x) \text {, a.e. in } \Omega\right\} \text { such that: } \\
& \qquad \int_{\Omega} \sum_{i=1}^{n} a_{i}(D u)\left(D_{i} v-D_{i} u\right) d x \geq \int_{\Omega} f(v-u) d x, \quad \forall v \in K_{g} \tag{2}
\end{align*}
$$

where $a(p): \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a strongly monotone operator of class $C^{2}$ (see (6)).
Let us note that it follows from classical results in the literature that there exists a unique solution to problem (2) (see [25]).

In the first result of the paper (Theorem 1) we show that, under a condition on the gradient constraint $g$, problem (2) is equivalent to the following double obstacle problem

Find $u \in K$ such that:

$$
\begin{equation*}
\int_{\Omega} \sum_{i=1}^{n} a_{i}(D u)\left(D_{i} v-D_{i} u\right) d x \geq \int_{\Omega} f(v-u) d x, \quad \forall v \in K, \tag{3}
\end{equation*}
$$

where $K=\left\{v \in H_{0}^{1,2}(\Omega): w_{1}(x) \leq v(x) \leq w_{2}(x)\right.$ a.e. in $\left.\Omega\right\}$, and

$$
w_{1}=\inf _{v \in K} v(x), \quad w_{2}=\sup _{v \in K} v(x) .
$$

From Theorem 5.1 in [26] (see also [20]), $w_{2} \in H^{1, \infty}(\Omega)$ is the viscosity solution to the Hamilton-Jacobi equation

$$
\left\{\begin{array}{cc}
|D u|=\sqrt{g(x)} & \text { a.e. in } \Omega  \tag{4}\\
u=0 & \text { on } \partial \Omega
\end{array}\right.
$$

and

$$
\begin{equation*}
w_{2}(x)=\inf _{x_{0} \in \partial \Omega} L\left(x, x_{0}\right) \tag{5}
\end{equation*}
$$

where

$$
\begin{gathered}
L\left(x, x_{0}\right)= \\
\inf \left\{\int_{0}^{T_{0}} \sqrt{g(\xi(s))} d s: \xi:\left[0, T_{0}\right] \rightarrow \bar{\Omega}, \xi(0)=x, \xi\left(T_{0}\right)=x_{0},\left|\xi^{\prime}(s)\right| \leq 1 \text { a.e. in }\left[0, T_{0}\right]\right\}
\end{gathered}
$$

$w_{1}$ can be calculated analagously.
It is important to note that the regularity of these obstacles follows from the theory of the viscosity solutions to Hamilton-Jacobi equations (see [17], p. 31), even if the solutions to the Hamilton-Jacobi equations are, in general, not smooth.

As already recalled, the problems are, generally, not equivalent, but a condition on the sign of the second derivatives of $g$ is required.

Before proving the equivalence, in Section 3 we achieve a regularity result for solutions to (3) (Theorem 8), that we need in the sequel.

Then, thanks to the equivalence, it is possible to prove that Lagrange multipliers exist in $L^{2}$ (Theorem 2).

Finally, an existence result of Lagrange multipliers as Radon measures holds, under less restrictive assumptions (Theorem 3).

The results are obtained following variational arguments and the strong duality theory.
Let us remark, that, during the past several decades, the variational methods have played a key role in solving many problems arising in nonlinear analysis and optimization theory such as differential hemivariational inequalities systems (see [27]), monotone bilevel equilibrium problems, generalized global fractional-order composite dynamical systems, generalized time-dependent hemivariational inequalities systems, optimal control of feedback control system, and so on.

Moreover, let us emphasize that real-life applications have been investigated on the basis of the theory of variational inequalities with operators of monotone type (see [28-30] for mathematical models describing flows of Bingham-type fluids and flows of an Oldroyd type by means of a variational inequality approach).

Finally, let us stress that the problem under consideration is strictly connected to the Monge-Kantorovich mass transfer problem. In particular, in [31] the authors study the integrability of the Lagrange multiplier, assuming that $f$ belongs to $L^{p}(\Omega)$ in the case of constant gradient constraint (see also [32] for variable constraint $g$ ). The Monge-Kantorovich mass transfer problem has applications in diverse fields such as economics, image processing, computer vision, transportation planning, and statistical physics. It provides a mathematical framework for studying the optimal flow of mass, resources, or information between different distributions or regions.

The paper is organized as follows: in Section 2 we state our main results of equivalence between the variational problems and existence of Lagrange multipliers, in Section 3 we provide a preliminary regularity result and some results of the theory of strong duality are recalled. In Section 4 we prove Theorem 1 and Section 5 is devoted to the proofs of Theorems 2 and 3. Finally, in Section 6 we provide our conclusions and suggest new problems that may be of interest for future research.

## 2. Results

The main results of the paper are presented in this section.
In what follows we assume that $\Omega$ is an open bounded convex subset of $\mathbb{R}^{n}$ and the boundary $\partial \Omega$ is of class $C^{2}$.

Moreover, the operator $a$ is of class $C^{2}$, with $a(0)=0$.
In the first two results, we assume that $a$ is a strongly monotone operator, that is, there exists $\lambda>0$, such that

$$
\begin{equation*}
(a(p)-a(q), p-q) \geq \lambda\|p-q\|^{2} \quad \forall p, q \in \mathbb{R}^{n}, p \neq q . \tag{6}
\end{equation*}
$$

Theorem 1. Assuming that a satisfies assumption (6), $f \equiv$ constant $>0$, and the following condition is fulfilled

$$
\begin{equation*}
-\sum_{i, j=1}^{n} \frac{\partial}{\partial x_{i}}\left(\frac{\partial a_{i}}{\partial p_{j}} \frac{\partial g}{\partial x_{j}}\right) \geq 0 \quad \text { in } \Omega \tag{7}
\end{equation*}
$$

then, the solution $u$ to problem (2) is also the solution to problem (3).
Moreover, the following coincidence of sets holds:

$$
P=\left\{x \in \Omega:|D u|^{2}=g(x)\right\}=I=\left\{x \in \Omega: u(x)=w_{1}(x) \text { or } u(x)=w_{2}(x)\right\} .
$$

Regarding the Lagrange multipliers, we prove the existence in two different cases. In the first one, the Lagrange multipliers are $L^{2}$ functions, whereas, in the second one, under less restrictive assumptions, they are measures.

Let us stress that the second result (Theorem 3) holds under assumption of strictly monotonicity on the operator $a$, namely

$$
\begin{equation*}
(a(P)-a(Q), P-Q)>0 \quad \forall P, Q \in \mathbb{R}^{n}, P \neq Q \tag{8}
\end{equation*}
$$

Theorem 2. Under the same assumptions as in Theorem 1, if $u \in K_{g} \cap W^{2, p}(\Omega)$ solves problem (2), then, there exists a Lagrange multiplier $v \in L^{2}(\Omega), v \geq 0$ a.e. in $\Omega$, that is

$$
\left\{\begin{array}{l}
v\left(\sum_{i=1}^{n}\left(D_{i} u\right)^{2}-g(x)\right)=0 \text { a.e. in } \Omega  \tag{9}\\
\sum_{i=1}^{n} \frac{\partial a_{i}(D u)}{\partial x_{i}}+f=v \text { a.e. in } \Omega .
\end{array}\right.
$$

Theorem 3. Assume that a satisfies assumption (8) and $f \in L^{p}(\Omega), p>1$. If $u \in K_{g}$ solves problem (2), then there exists a Lagrange multiplier $\mu^{*} \in\left(L^{\infty}(\Omega)\right)^{*}$, that is

$$
\left\{\begin{array}{l}
\left\langle\mu^{*}, y\right\rangle \geq 0 \quad \forall y \in L^{\infty}(\Omega), y \geq 0 \quad \text { a.e. in } \Omega  \tag{10}\\
\left.\left\langle\mu^{*}, \sum_{i=1}^{n}\left(D_{i} u\right)^{2}-g(x)\right)\right\rangle=0 ; \\
\int_{\Omega}\left\{\sum_{i=1}^{n} a_{i}(D u) \frac{\partial \varphi}{\partial x_{i}}-f \varphi\right\} d x=\left\langle\mu^{*},-2 \sum_{i=1}^{n} \frac{\partial u}{\partial x_{i}} \frac{\partial \varphi}{\partial x_{i}}\right\rangle \quad \forall \varphi \in H_{0}^{1, \infty}(\Omega) .
\end{array}\right.
$$

## 3. Preliminary Results

This section is devoted to some preliminary results that we need to prove our theorems.
In particular, first we recall the strong duality theory and, then, we prove a regularity result for the solution to the double obstacle problem (3) that we need to apply in Section 4 a maximum principle.

For the sake of clarity, here we provide the main results of classical strong duality theory and a new strong duality theory, obtained using new separation theorems based on the notion of quasi-relative interior.

For the classical results of strong duality theory we refer to ([33], Theorems 6.7 and 6.11).
It is important to note that strong duality has important implications in optimization theory. It allows us to obtain lower bounds on the optimal value of the primal problem by solving the dual problem. It also provides a way to assess optimality and obtain dual solutions that can provide additional information about the primal problem, such as shadow prices or sensitivity analysis.

The framework, in which the classical theory works, is the following one: $X$ is a real linear space and $S \subset X$ is a nonempty subset; $(Y,\|\cdot\|)$ is a partially ordered real normed space with ordering cone $C$, and $C^{*}=\left\{\lambda \in Y^{*}:\langle\lambda, y\rangle \geq 0 \forall y \in C\right\}$ is the dual cone of $C$, whereas $Y^{*}$ is the topological dual of $Y$. Moreover, $F: S \rightarrow \mathbb{R}$ is a given objective
functional, $G: S \rightarrow Y$ is a given constraint mapping and the constraint set is given as $\mathbb{K}:=\{v \in S: G(v) \in-C\}$.

We consider the primal problem

$$
\begin{equation*}
\min _{\substack{G(v) \in-C \\ v \in S}} F(v) \tag{11}
\end{equation*}
$$

and the dual problem

$$
\begin{equation*}
\max _{\lambda \in C^{*}} \inf _{v \in S}[F(v)+\lambda(G(v))] \tag{12}
\end{equation*}
$$

where $\lambda$ is the Lagrange multiplier associated with the sign constraints.
As is well known (see [33]), the weak duality always holds, namely,

$$
\begin{equation*}
\max _{\lambda \in C^{*}} \inf _{v \in S}[F(v)+\lambda(G(v))] \leq \min _{G(v) \in-C}^{v \in S}(F(v) \tag{13}
\end{equation*}
$$

Moerover, if problem (11) is solvable and in (13) the equality holds, the strong duality between the primal problem (11) and the dual problem (12) holds.

Theorem 4 (classical strong duality property [33]). Assume that the composite mapping $(F, G)$ : $S \rightarrow \mathbb{R} \times Y$ is convex-like with respect to product cone $\mathbb{R}_{+} \times C$ in $\mathbb{R} \times Y, \mathbb{K}$ is nonempty and the ordering cone $C$ has a nonempty interior int $(C)$. If the primal problem (11) is solvable and the generalized Slater condition is satisfied, namely there is a vector $\bar{v} \in S$ with $G(\bar{v}) \in-\operatorname{int}(C)$, then the dual problem (12) is also solvable and the extremal values of the two problems are equal. Moreover, if $u$ is the optimal solution to problem (11) and $\bar{v} \in C^{*}$ is a solution to problem (12), it follows that

$$
\begin{equation*}
\bar{v}(G(u))=0 . \tag{14}
\end{equation*}
$$

Moreover, if

$$
\mathcal{L}(v, v)=F(v)+v(G(v))
$$

is the Lagrange functional, then the following relationship with the saddle points of $\mathcal{L}(v, v)$ holds.

Theorem 5 (see [33]). Under the same assumptions as in Theorem 4, if the ordering cone $C$ is closed, then a point $(u, \bar{v}) \in S \times C^{*}$ is a saddle point of the Lagrange functional $\mathcal{L}(v, v)$, namely

$$
\mathcal{L}(u, v) \leq \mathcal{L}(u, \bar{v}) \leq \mathcal{L}(v, \bar{v}), \forall v \in S, \forall v \in C^{*},
$$

if and only if $u$ is a solution to the primal problem (11), $\bar{v}$ is a solution to the dual problem (12) and the extremal values of the two problems are equal.

Let us stress that we apply classical strong duality theory to prove Theorem 3, whereas we need a new theory (see [9]) to obtain the other results. Indeed, in our framework, as in many applications in infinite dimensional settings, the classical theory does not work, since the assumption of nonemptiness of the ordering cone is not fulfilled.

Here, we recall the new strong duality theory in its complete version, namely in the case of inequality and equality constraints.

The assumptions read as follows:
Let $\left(X,\|\cdot\|_{X}\right),\left(Y,\|\cdot\|_{Y}\right),\left(Z,\|\cdot\|_{Z}\right)$ be real normed spaces with $Y^{*}, Z^{*}$ topological dual of $Y$ and $Z$, respectively; $Y$ is partially ordered by a convex cone $C, C^{*}=\left\{\mu \in Y^{*}\right.$ : $\langle\mu, y\rangle \geq 0 \forall y \in C\}$ is the dual cone of $C$. $S$ is a nonempty subset of $X$, and $F: S \rightarrow \mathbb{R}$, $G: S \rightarrow Y, H: S \rightarrow Z$ are three functions.

Moreover, we define the feasible set

$$
\mathbb{K}=\left\{v \in S: G(v) \in-C, H(v)=\theta_{Z}\right\}
$$

We recall the definition of tangent cone to $S^{*} \subset X$ at a point $v \in X$ :

$$
T_{S^{*}}(v):=\left\{l \in X: l=\lim _{n} \mu_{n}\left(v_{n}-v\right), \mu_{n}>0, v_{n} \in S^{*} \forall n \in N, \lim _{n} v_{n}=v\right\}
$$

We introduce the following constraint qualification assumption: we say that Assumption $S$ is satisfied at a point $v_{0} \in \mathbb{K}$ if

$$
\begin{equation*}
T_{\widetilde{N}}\left(0, \theta_{Y}, \theta_{Z}\right) \cap(]-\infty, 0\left[\times\left\{\theta_{Y}\right\} \times\left\{\theta_{Z}\right\}\right)=\varnothing \tag{15}
\end{equation*}
$$

where

$$
\widetilde{N}=\left\{\left(F(v)-F\left(v_{0}\right)+\alpha, G(v)+w, H(v)\right): v \in S \backslash \mathbb{K}, \alpha \geq 0, w \in C\right\}
$$

Under Assumption S the following strong duality property holds (see [9]).
Theorem 6. Let us assume that $F$ and $G$ are convex functions, $H$ is an affine-linear mapping and $v_{0} \in \mathbb{K}$ is a solution to the primal problem

$$
\begin{equation*}
\min _{v \in \mathbb{K}} F(v) . \tag{16}
\end{equation*}
$$

Then, if Assumption $S$ is fulfilled at $v_{0}$, the dual problem

$$
\begin{equation*}
\max _{\substack{\lambda \in C^{*} \\ \mu \in Z^{*}}} \inf _{v \in S}\{F(v)+\langle\lambda, G(v)\rangle+\langle\mu, H(v)\rangle\} \tag{17}
\end{equation*}
$$

is also solvable and the extreme values of the primal problem and of the dual problem coincide. Moreover, if $\left(v_{0}, \lambda^{*}, \mu^{*}\right) \in \mathbb{K} \times C^{*} \times Z^{*}$ solves problem (17), then $\left\langle\lambda^{*}, G\left(v_{0}\right)\right\rangle=0$.

Moreover, if

$$
\mathcal{L}(v, \lambda, \mu)=F(v)+\langle\lambda, G(v)\rangle+\langle\mu, H(v)\rangle
$$

is the Lagrange functional, then the following result on the saddle points of the Lagrange functional holds.

Theorem 7 ([9]). Under the same assumptions as in Theorem 6, $v_{0} \in \mathbb{K}$ solves problem (16) if and only if there exist $\lambda^{*} \in C^{*}$ and $\mu^{*} \in Z^{*}$ such that $\left(x_{0}, \lambda^{*}, \mu^{*}\right)$ is a saddle point of the Lagrange functional, namely

$$
\mathcal{L}\left(v_{0}, \lambda, \mu\right) \leq \mathcal{L}\left(v_{0}, \lambda^{*}, \mu^{*}\right) \leq \mathcal{L}\left(v, \lambda^{*}, \mu^{*}\right), \quad \forall v \in S, \lambda \in C^{*}, \mu \in Z^{*} .
$$

Now, we prove the following regularity result, that we will use in Section 4.

Theorem 8. Let the assumptions of Theorem 1 be satisfied and $u$ be the solution to problem (3). Then, $u \in W^{2, p}(\Omega)$. In particular, if $p>n, D u \in C^{0, \alpha}(\Omega)$.

Proof. The first goal is an estimate for

$$
|u|_{1}=\sup \left\{\frac{|u(x)-u(y)|}{|x-y|}: x, y, \in \bar{\Omega}, x \neq y\right\},
$$

obtained using similar arguments as in [34].
Let $u$ be the solution to (3), we set $\tilde{u}$ the extension by zero of $u$ to $\mathbb{R}^{n}$ and

$$
\begin{equation*}
u^{h}(x)=\max \{\tilde{u}(x+h)-\tilde{u}(x)-M|h|, 0\} \quad \forall x, h \in R^{n}, \tag{18}
\end{equation*}
$$

where $M=\max \left\{\left|w_{1}\right|_{1},\left|w_{2}\right|_{1}\right\}$.

## Defining

$$
\begin{gathered}
u_{1}(x)=\max \{\tilde{u}(x), \tilde{u}(x+h)-M|h|\}=\tilde{u}(x)+u^{h}(x), \\
u_{2}(x)=\min \{\tilde{u}(x), \tilde{u}(x-h)+M|h|\}=\tilde{u}(x)-u^{h}(x-h) .
\end{gathered}
$$

as in [19,23], we have $u_{1 / \Omega}, u_{2 / \Omega} \in K$ and

$$
\tilde{w}_{1}(x) \leq u_{2}(x) \leq u_{1}(x) \leq \tilde{w}_{2}(x) \text { a.e. in } \Omega .
$$

Following the same arguments as in [12], we get

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} \sum_{i=1}^{n}\left(a_{i}(D \tilde{u}(x+h))-a_{i}(D \tilde{u}(x))\right) D_{i} u^{h}(x) d x \leq 0 . \tag{19}
\end{equation*}
$$

Setting $X_{h}^{+}=\left\{x \in R^{n}: \tilde{u}(x+h)-\tilde{u}(x)-M|h| \geq 0\right\}$, from (18) and (19) it follows that

$$
\begin{equation*}
\int_{X_{h}^{+}} \sum_{i=1}^{n}\left(a_{i}(D \tilde{u}(x+h))-a_{i}(D \tilde{u}(x))\right)\left(D_{i} \tilde{u}(x+h)-D_{i} \tilde{u}(x)\right) d x \leq 0 . \tag{20}
\end{equation*}
$$

Thanks to strong monotonicity assumption (6) and to inequality (20), we may conclude that $u^{h}=0$ in $X_{h}^{+}$and then

$$
\tilde{u}(x+h)-\tilde{u}(x)-M|h| \leq 0 \quad \forall x, h \in \mathbb{R}^{n}
$$

namely,

$$
|u|_{1} \leq M
$$

and

$$
\begin{equation*}
|D u| \leq M \quad \text { a.e. in } \Omega . \tag{21}
\end{equation*}
$$

To conclude, we consider the following elastic-plastic torsion problem
Find $w \in K_{M}=K \cap\left\{v \in H_{0}^{1}:|D v| \leq M\right.$ a.e. in $\left.\Omega\right\}$ such that:

$$
\begin{equation*}
\int_{\Omega} \sum_{i=1}^{n} a_{i}(D w)\left(D_{i} v-D_{i} w\right) d x \geq \int_{\Omega} f(v-w) d x, \quad \forall v \in K_{M} \tag{22}
\end{equation*}
$$

Since the feasible set $K_{M}$ is a bounded, closed, and convex set, from classical results (see [35]), the unique solution $u \in K_{M}$ to the variational inequality (22) belongs to $W^{2, p}(\Omega)$. Then, the thesis is achieved.

## 4. The Equivalence of the Two Variational Problems

Now, we may prove Theorem 1.
Obviously,

$$
\begin{equation*}
K_{g} \subseteq K . \tag{23}
\end{equation*}
$$

Then, to prove the equivalence of the two problems, we have to show that if $u \in K$ is the solution to (3), then $u$ belongs to $K_{g}$.

To this aim, setting

$$
F(v)=\int_{\Omega}\left\{\sum_{i=1}^{n} a_{i}(D u)\left(D_{i} v-D_{i} u\right)-f(v-u)\right\} d x
$$

we note that problem (3) may be rewritten as the optimization problem

$$
\begin{equation*}
\min _{v \in K} F(v), \tag{24}
\end{equation*}
$$

which satisfies Assumption S.

Indeed, if we set

$$
\begin{gathered}
X=S=L^{2}(\Omega), Y=L^{2}(\Omega) \times L^{2}(\Omega), \\
C=C^{*}=\left\{(a(x), b(x)) \in L^{2}(\Omega) \times L^{2}(\Omega): a(x), b(x) \geq 0 \text { a.e. in } \Omega\right\}, \\
G(v)=\left(G_{1}(v), G_{2}(v)\right)=\left(w_{1}-v, v-w_{2}\right),
\end{gathered}
$$

we have

$$
\tilde{N}=\left\{\left(F(v)+\alpha, w_{1}-v+a, v-w_{2}+b\right), v \in L^{2} \backslash K, \alpha \geq 0, y=(a, b) \in C\right\}
$$

Following similar arguments as in $[23,36]$ we may show that, if

$$
\left(l, \theta_{L^{2}(\Omega)}, \theta_{L^{2}(\Omega)}\right)=\lim _{n}\left[\mu_{n}\left(F\left(v_{n}\right)+\alpha_{n}, w_{1}-v_{n}+a_{n}, v_{n}-w_{2}+b_{n}\right)\right]
$$

with $\mu_{n}>0, \lim _{n}\left(F\left(v_{n}\right)+\alpha_{n}\right)=0, \alpha_{n} \geq 0, v_{n} \in L^{2}(\Omega) \backslash K, \lim _{n} \mu_{n}\left(w_{1}-v_{n}+a_{n}\right)=$ $\theta_{L^{2}(\Omega)}, \lim _{n} \mu_{n}\left(v_{n}-w_{2}+b_{n}\right)=\theta_{L^{2}(\Omega)}, y_{n}=\left(a_{n}, b_{n}\right) \in C$, then

$$
l \geq 0
$$

namely, Assumption $S$ is fulfilled at the solution to problem (24).
Then, if we consider the Lagrange functional

$$
\begin{gather*}
\mathcal{L}(v, \lambda, \mu)=  \tag{25}\\
=\int_{\Omega}\left(-\sum_{i=1}^{n} \frac{\partial a_{i}(D u)}{\partial x_{i}}-f\right)(v-u) d x+\int_{\Omega} \lambda\left(w_{1}(x)-v(x)\right) d x+\int_{\Omega} \mu\left(v(x)-w_{2}(x)\right) d x
\end{gather*}
$$

thanks to Theorem 7, there exists a saddle point $\left(\lambda^{*}, \mu^{*}\right) \in C$, namely,

$$
\begin{equation*}
\mathcal{L}(u, \lambda, \mu) \leq \mathcal{L}\left(u, \lambda^{*}, \mu^{*}\right) \leq \mathcal{L}\left(v, \lambda^{*}, \mu^{*}\right) \quad \forall v \in L^{2}(\Omega), \forall(\lambda, \mu) \in C \tag{26}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{\Omega} \lambda^{*}\left(w_{1}(x)-u(x)\right) d x=0, \quad \int_{\Omega} \mu^{*}\left(u(x)-w_{2}(x)\right) d x=0 \tag{27}
\end{equation*}
$$

that is,

$$
\begin{equation*}
\lambda^{*}\left(w_{1}(x)-u(x)\right)=0, \quad \mu^{*}\left(u(x)-w_{2}(x)\right)=0, \quad \text { a.e. in } \Omega . \tag{28}
\end{equation*}
$$

Using variational arguments (see [12]), it follows

$$
\begin{equation*}
-\sum_{i=1}^{n} \frac{\partial a_{i}(D u)}{\partial x_{i}}-f-\lambda^{*}+\mu^{*}=0 \quad \text { a.e. in } \Omega . \tag{29}
\end{equation*}
$$

Now, we consider the coincidence set $I=\left\{x \in \Omega: u(x)=w_{1}(x)\right.$ or $\left.u(x)=w_{2}(x)\right\}$ and the non-coincidence set $N=\left\{x \in \Omega: w_{1}(x)<u(x)<w_{2}(x)\right\}$.

From [26], Theorem 5.1, we have that $\left|D w_{1}(x)\right|=\left|D w_{2}(x)\right|=\sqrt{g(x)}$ a.e. in $\Omega$, then

$$
\begin{equation*}
|D u|=\sqrt{g(x)} \text { in } I \tag{30}
\end{equation*}
$$

Moreover, from (28) and (29) it follows that $\lambda^{*}=\mu^{*}=0$ a.e. in $N$ and

$$
\begin{equation*}
-\sum_{i=1}^{n} \frac{\partial a_{i}(D u)}{\partial x_{i}}=f \quad \text { a.e. in } N . \tag{31}
\end{equation*}
$$

Thanks to the regularity of $u$, stated in Theorem 8 , and since $f$ is a constant function, we follow the same steps used in [35], Lemma III.10. We differentiate (31) with respect to $x_{k}$, multiply it by $\frac{\partial u}{\partial x_{k}}$ and sum it with respect to $k$. Then, it follows that

$$
\begin{equation*}
\sum_{i, j, k=1}^{n} \frac{\partial}{\partial x_{k}}\left(\frac{\partial a_{i}(D u)}{\partial p_{j}}\right) \frac{\partial^{2} u}{\partial x_{j} \partial x_{i}} \frac{\partial u}{\partial x_{k}}+\sum_{i, j, k=1}^{n} \frac{\partial a_{i}(D u)}{\partial p_{j}} \frac{\partial^{3} u}{\partial x_{j} \partial x_{i} \partial x_{k}} \frac{\partial u}{\partial x_{k}}=0 . \tag{32}
\end{equation*}
$$

It follows that

$$
\begin{gather*}
\frac{1}{2} \frac{\partial}{\partial x_{i}}\left[\sum_{j=1}^{n} \frac{\partial a_{i}(D u)}{\partial p_{j}} \frac{\partial}{\partial x_{j}}\left(|D u|^{2}-g(x)\right)\right]=\sum_{i, j, k=1}^{n} \frac{\partial}{\partial x_{i}}\left(\frac{\partial a_{i}(D u)}{\partial p_{j}}\right) \frac{\partial u}{\partial x_{k}} \frac{\partial^{2} u}{\partial x_{k} \partial x_{j}} \\
+\sum_{i, j, k=1}^{n} \frac{\partial a_{i}(D u)}{\partial p_{j}} \frac{\partial^{2} u}{\partial x_{k} \partial x_{j}} \frac{\partial^{2} u}{\partial x_{k} \partial x_{i}}+\sum_{i, j, k=1}^{n} \frac{\partial a_{i}(D u)}{\partial p_{j}} \frac{\partial u}{\partial x_{k}} \frac{\partial^{3} u}{\partial x_{j} \partial x_{i} \partial x_{k}}  \tag{33}\\
-\frac{1}{2} \sum_{i, j=1}^{n} \frac{\partial}{\partial x_{i}}\left(\frac{\partial a_{i}(D u)}{\partial p_{j}} \frac{\partial g(x)}{\partial x_{j}}\right) .
\end{gather*}
$$

From assumptions (6) and (7), we have

$$
\begin{array}{r}
\frac{1}{2} \frac{\partial}{\partial x_{i}}\left[\sum_{j=1}^{n} \frac{\partial a_{i}(D u)}{\partial p_{j}} \frac{\partial}{\partial x_{j}}\left(|D u|^{2}-g(x)\right)\right] \\
\geq \sum_{i, j, k=1}^{n}\left[\frac{\partial}{\partial x_{i}}\left(\frac{\partial a_{i}(D u)}{\partial p_{j}}\right) \frac{\partial^{2} u}{\partial x_{k} \partial x_{j}}-\frac{\partial}{\partial x_{k}}\left(\frac{\partial a_{i}(D u)}{\partial p_{j}}\right) \frac{\partial^{2} u}{\partial x_{j} \partial x_{i}}\right] \frac{\partial u}{\partial x_{k}}  \tag{34}\\
=\sum_{i, j, k=1}^{n}\left[\sum_{l=1}^{n} \frac{\partial^{2} a_{i}(D u)}{\partial p_{j} \partial p_{l}} \frac{\partial^{2} u}{\partial x_{l} \partial x_{i}} \frac{\partial^{2} u}{\partial x_{k} \partial x_{j}}-\sum_{l=1}^{n} \frac{\partial^{2} a_{i}(D u)}{\partial p_{j} \partial p_{l}} \frac{\partial^{2} u}{\partial x_{k} \partial x_{l}} \frac{\partial^{2} u}{\partial x_{j} \partial x_{i}}\right] \frac{\partial u}{\partial x_{k}}=0 .
\end{array}
$$

Finally, since the coefficients are bounded, $N$ is an open set, applying the maximum principle to the operator

$$
-\mathcal{A}(\varphi)=-\frac{1}{2} \sum_{i, j=1}^{n} \frac{\partial}{\partial x_{i}}\left(\frac{\partial a_{i}}{\partial p_{j}} \frac{\partial \varphi}{\partial x_{j}}\right)
$$

acting on $|D u|^{2}-g(x)$ on $N$, we have

$$
\begin{equation*}
|D u(x)|<\sqrt{g(x)} \quad \text { a.e. } \in N . \tag{35}
\end{equation*}
$$

From (30) and (35) it follows that, if $u \in K$ is a solution to (3), then

$$
\begin{equation*}
|D u(x)| \leq \sqrt{g(x)} \quad \text { a.e. } \in \Omega \tag{36}
\end{equation*}
$$

Taking into account the uniqueness of the solution, we may conclude that the solution to (3) is also the solution to (2) and Theorem 1 is proved.

Finally, the following interesting coincidence of sets follows from (30) and (35)

$$
E=\{x \in \Omega:|D u|<\sqrt{g(x)}\}=N .
$$

## 5. Lagrange Multipliers

In this section we provide the proofs of the existence of Lagrange multipliers.
A first result, the existence of $L^{2}$ Lagrange multipliers, holds under the assumption $f \equiv$ constant $>0$ and $a$ strongly monotone operator. It follows from (28) and (29) as in the proof of Theorem 1.

The second result holds assuming that $f \in L^{p}(\Omega), p>1$, and the operator $a$ is strictly monotone. In this case the Lagrange multipliers exist in the dual of $L^{\infty}$.

Indeed, we set

$$
X=S=W_{0}^{1, \infty}(\Omega) ; \quad C=\left\{v \in L^{\infty}(\Omega): v(x) \geq 0 \text { a.e. in } \Omega\right\} .
$$

In this case $C$ has a nonempty interior, then we may apply the classical strong duality theory (see [33]).

We may rewrite problem (2) as
Find $u \in K_{1}=\left\{v \in H_{0}^{1, \infty}(\Omega): \sum_{i=1}^{n}\left(\frac{\partial v}{\partial x_{i}}\right)^{2} \leq g(x)\right.$, a.e. on $\left.\Omega\right\}$ such that:

$$
\begin{equation*}
\int_{\Omega}\left\{\sum_{i=1}^{n} a_{i}(D u)\left(D_{i} v-D_{i} u\right)-f(v-u)\right\} d x \geq 0, \quad \forall v \in K_{1} \tag{37}
\end{equation*}
$$

Following the same steps as in [19], we may prove that $C$ is closed and the generalized Slater condition is verified. Moreover, since $F$ and $G$ are convex, then the composite mapping $(F, G)$ is convex-like, namely all the assumptions of Theorems 6.7 and 6.11 in [33] are fulfilled.

Then, it follows that there exists $\mu^{*} \in C^{*}$ solution to the dual problem

$$
\begin{equation*}
\max _{\mu \in C^{*}} \inf _{v \in S}[F(v)+\langle\mu, G(v)\rangle], \tag{38}
\end{equation*}
$$

with

$$
\begin{equation*}
F(v)=\int_{\Omega}\left\{\sum_{i=1}^{n} a_{i}(D u)\left(D_{i} v-D_{i} u\right)-f(v-u)\right\} d x \tag{39}
\end{equation*}
$$

and

$$
G(v)=|D v|^{2}-g(x) .
$$

Moreover, $\left(u, \mu^{*}\right)$ is a saddle point of the Lagrange functional

$$
\mathcal{L}(v, \mu)=F(v)+\langle\mu, G(v)\rangle, \forall v \in H_{0}^{1, \infty}(\Omega), \forall \mu \in C^{*},
$$

that is

$$
\begin{equation*}
\mathcal{L}(u, \mu) \leq \mathcal{L}\left(u, \mu^{*}\right) \leq \mathcal{L}\left(v, \mu^{*}\right), \forall v \in H_{0}^{1, \infty}(\Omega), \forall \mu \in C^{*} \tag{40}
\end{equation*}
$$

Using variational arguments as in [19], we obtain that $\mu^{*} \in\left(L^{\infty}(\Omega)\right)^{*}$ satisfies conditions (10).

## 6. Discussions

The paper adds to the already existing literature on nonconstant gradient constrained problem further results related to the relationship with double obstacle problem and the existence of Lagrange multipliers.

In particular, in the paper we focused on the nonconstant gradient constraint $|D u| \leq g(x)$ associated with a nonlinear monotone operator $a(D u)$.

The existence of Lagrange multipliers as Lebesgue functions is guaranteed in the case $f \equiv$ constant $>0$ and strong monotonicity assumption on the operator, whereas the Lagrange multipliers exist as Radon measure in the case $f \in L^{p}, p>1$, and strict monotonicity assumption is required.

In the future, several studies could be carried out in several directions in this framework. For example it will be interesting to consider a regular, nonconstant, free term $f$, or studying the problem associated with different nonlinear operators. Moreover, the properties of the Lagrange multiplier may be investigated. Finally, one could analyze the natural parabolic counterpart.

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# New Class of K-G-Type Symmetric Second Order Vector Optimization Problem 

Chetan Swarup ${ }^{1(D)}$, Ramesh Kumar ${ }^{2}$, Ramu Dubey ${ }^{2, *}$ and Dowlath Fathima ${ }^{3}$ (D)<br>1 Department of Basic Science, College of Science and Theoretical Studies, Saudi Electronic University, Riyadh-Male Campus, Riyadh 13316, Saudi Arabia; c.swarup@seu.edu.sa<br>${ }^{2}$ Department of Mathematics, J.C. Bose University of Science and Technology, YMCA, Faridabad 121 006, India; rameshgoswami2619@gmail.com<br>3 Basic Sciences Department, College of Science and Theoretical Studies, Saudi Electronic University, Jeddah 23442, Saudi Arabia; d.fathima@seu.edu.sa<br>* Correspondence: rdubeyjiya@gmail.com

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#### Abstract

In this paper, we present meanings of $K-G_{f}$-bonvexity $/ K-G_{f}$-pseudobonvexity and their generalization between the above-notice functions. We also construct various concrete non-trivial examples for existing these types of functions. We formulate $K-G_{f}$-Wolfe type multiobjective secondorder symmetric duality model with cone objective as well as cone constraints and duality theorems have been established under these aforesaid conditions. Further, we have validates the weak duality theorem under those assumptions. Our results are more generalized than previous known results in the literature.


Keywords: $K$ - $G_{f}$-pseudobonvexity; second-order; $K$ - $G_{f}$-Wolfe type; efficient solution; multiobjective programming; arbitrary cones; strong duality; generalized assumptions

MSC: 90C26; 90C30; 90C32; 90C46

## 1. Introduction

The field of optimization theory has progressed far beyond anyone's expectations. Due to its wide variety of uses, it has made its way into all disciplines of science and engineering. When approximations are utilized, one of the most important practical applications of duality is that it provides bounds on the value of the objective functions because there are more factors involved, second-order duality has a greater computational benefit than firstorder duality. For intriguing applications and breakthroughs in multiobjective optimization, we refer to [1] , and the references cited therein. Dorn [2] presented the primary symmetric duality definition for quadratic programming in 1965. Dantzig et al. [3] and Mond [4] proposed a pair of symmetric dual Duality plays a vital role in investigating nonlinear programming problem solutions. Several writers have proposed several duality models, such as Wolfe dual [5] and Mond-Weir dual [6]. Nanda and Das [7] introduced four different forms of duality models for the nonlinear programming problem with cone constraints. The work of Bazaraa and Goode [8] and Hanson and Mond [9] inspired these findings.

Mangasian [10] established the duality theorem in the context of a second-order dual problem in nonlinear programming, where none of the constraints imposed convexity restrictions on all functions. Mond [11] introduced second-order symmetric dual models and established second-order symmetric duality theorems under second-order convexity conditions for the first time. In mathematical programming, Hasnson [12] defined the second-order invexity of a differentiable function and studied it. In 1999, Mishra [13] proposed a pair of second-order vector symmetric dual multiobjective models for arbitrary cones based on the Wolfe and Mond-Weir types. In addition 2006, ref. [14] a couple of Mond-Weir type second-order symmetric duality multiobjective calculations for cone
second-order pseudoinvex and emphatically cone second-order pseudoinvex algorithm were presented. A couple of Mond-Weir type second-order symmetric dual multiobjective projects over discretion cones is created under pseudoinvexity $/ K\ulcorner F$-convexity assumptions by Gulati [15], which is as:
Primal(MP):

$$
\begin{aligned}
& \begin{array}{l}
\text { K-minimize } \\
\text { subject to } \\
\\
-\quad\left(\nabla_{\kappa}\left(\lambda^{T} \psi\right)(\iota, \kappa)+\nabla_{\kappa \kappa}\left(w^{T} \phi\right)(\iota, \kappa) p\right) \in C_{2}^{*}, \\
\\
\kappa^{T}\left(\nabla_{\kappa}\left(\lambda^{T} \psi\right)(\iota, \kappa)+\nabla_{\kappa \kappa}\left(w^{T} \phi\right)(\iota, \kappa) p\right) \geqq 0, \\
\lambda \in \operatorname{int}^{*}, \quad \iota \in C_{1}
\end{array}
\end{aligned}
$$

## Dual(MD):

$$
\begin{aligned}
& \begin{array}{l}
\text { K-maximize } \\
\text { subject to } \\
\\
\left(\nabla_{\iota}\left(\lambda^{T} \psi\right)(\mu, v)+\nabla_{l u}\left(w^{T} \phi\right)(\mu, v) p\right) \in C_{1}^{*}, \\
\\
\mu^{T}\left(\nabla_{\iota}\left(\lambda^{T} \psi\right)(\mu, v)+\nabla_{u}\left(w^{T} \phi\right)(\mu, v) r\right) \leqq 0, \\
\\
\lambda \in \operatorname{int}^{*}, \quad \iota \in C_{2},
\end{array}
\end{aligned}
$$

where,
(i) $R_{1} \subseteq \mathbb{R}^{n}, R_{2} \subseteq \mathbb{R}^{m}$ are open sets,
(ii) $\psi, \phi: R_{1} \times R_{2} \rightarrow \mathbb{R}^{k}$ is a twice differentiable function of $\iota$ and $\kappa$, is a differentiable function of $\iota$ and $\kappa$,
(iii) $\lambda \in \mathbb{R}^{k}, w \in \mathbb{R}^{q}, p \in \mathbb{R}^{m}$ and $r \in \mathbb{R}^{n}$,
(iv) for $i=1,2, C_{i} \subset S_{i}$ is a closed convex cone with non-empty interior and $C_{i}^{*}$ is its positive polar cone.
Aside from them, a number of other researchers are working in this field. For additional information, see [16-20].
In this paper be start by defining in section $2, K-G_{f}$-bonvexity as well as pseduobonvexity and construct non-trivial numerical examples for clear understanding the concept introduced by authors. We identify several examples lying exclusively K -$G_{f}$-bonvex and not in the class of $K$-invex function with respect to same $\eta$ already exist in the literature. We illustrate an example which is $K-G_{f}$-pseudobonvex but not $K-G_{f}$-bonvex with respect to same $\eta$. In the next section, we formulate a new pair of multiobjective symmetric second order $K-G_{f}$-primal-dual models over arbitrary cone and drive duality results under $K-G_{f}$-bonvex as well as $K-G_{f}$-pseudobonvex assumptions. We, also construct a non-trivial example for validate the weak duality theorem presented in the paper. we also introduced geometry figure for clear understanding the concept through figure.

## 2. Preliminaries and Definitions

In this paper, we used $\mathbb{R}^{n}$ for $n$-dimensional Euclidean space and $\mathbb{R}_{+}^{n}$ for semi-positive orthant. Also, here $C_{1}$ and $C_{2}$ used for closed convex cone $\mathbb{R}^{n}$ and $\mathbb{R}^{m}$ respectively, with non-void interiors. For a real-valued twice differentiable function $g(\varphi, \vartheta)$ described on an open set in $\mathbb{R}^{n} \times \mathbb{R}^{m}$, indicate by $\nabla_{\varphi} g(\bar{\varphi}, \bar{\vartheta})$ the gradient vector of $g$ with respect to $a$ at $(\bar{\varphi}, \bar{\vartheta}), \nabla_{\varphi \varphi} g(\bar{\varphi}, \bar{\vartheta})$ the Hessian matrix with respect to $\varphi$ an at $(\bar{\varphi}, \bar{\vartheta})$.
Throughout the paper $\tilde{N}=\{1,2, \ldots, k\}, \tilde{O}=\{1,2, \ldots, m\}$.
A differentiable function $f: X \times Y \rightarrow R^{k}, \eta_{1}: X \times Y \rightarrow R^{k}, \eta_{2}: X \times Y \rightarrow R^{k}$,
$G_{f}=\left(G_{f_{1}}, G_{f_{2}}, \ldots, G_{f_{k}}\right): R \rightarrow R^{k}, G_{f_{i}}: I_{f_{i}}(X) \rightarrow R$ is range $f_{i}$ for $i=\tilde{N}$. Also, $K$ is used for pointed convex cone with non-void interiors in $\mathbb{R}^{k}$, for $\vartheta, z \in \mathbb{R}^{k}$ and we specify cone orders with respect to $K$ as follows:

$$
\vartheta \leqq z \Longleftrightarrow z-\vartheta \in K ; \vartheta \leq z \Longleftrightarrow z-\vartheta \in K \backslash\{0\} ; \vartheta<z \Longleftrightarrow z-\vartheta \in \operatorname{int} K .
$$

Let $f: X \rightarrow \mathbb{R}^{k}$ be a differentiable function defined on open set $\phi \neq X \subseteq \mathbb{R}^{n}$ and $I_{f_{i}}(X), i \in \tilde{N}$ be the range of $f_{i}$.

Consider the following multiobjective programming problem with cone objective as well as constraints as :
(MP)

$$
\begin{aligned}
& K-\min f(\varphi) \\
& \text { subject to }
\end{aligned}
$$

$$
\varphi \in X^{0}=\{\varphi \in S: g(\varphi) \in Q\}
$$

where $S \subseteq \mathbb{R}^{n}, f: S \rightarrow \mathbb{R}^{k}, g: S \rightarrow \mathbb{R}^{m}$. $Q$ is a closed convex cone with a non-empty interior in $\mathbb{R}^{m}$.

Definition 1 ([21]). $\bar{\varphi} \in X^{0}$ is a weak efficient solution of (MP), $\nexists \varphi \in X$ such that

$$
f(\bar{\varphi})-f(\varphi) \in \operatorname{int} K .
$$

Definition 2 ([21]). $\bar{\varphi} \in X^{0}$ is an efficient solution of (MP), $\exists \varphi \in X$ such that

$$
f(\bar{\varphi})-f(\varphi) \in K \backslash\{0\} .
$$

Now, we consider the following multiobjective programming with cone objective and cone constraints as:
(GMP)

$$
\begin{aligned}
& K-\min G_{f}(f(z)) \\
& \text { subject to } z \in Z^{0}=\left\{z \in S:-G_{g}(g(z)) \in Q\right\} .
\end{aligned}
$$

Definition 3 ([21]). $\bar{z} \in Z^{0}$ is a weak efficient solution of (GMP), $\nexists z \in Z^{0}$ s.t.

$$
G_{f}(f(\bar{z}))-G_{f}(f(z)) \in \operatorname{int} K
$$

Definition 4 ([21]). $\bar{z} \in Z^{0}$ is a efficient solution of $(G M P), \nexists z \in Z^{0}$ s.t. $G_{f}(f(\bar{z}))-$ $G_{f}(f(z)) \in K \backslash\{0\}$.

Definition 5 ([21]). The positive polar cone $C_{i}^{*}$ of $C_{i}(i=1,2)$ is defined as $C_{i}^{*}=\left\{z: \varphi^{T} z \geqq\right.$ $\left.0, \forall \varphi \in C_{1}\right\}$.
Suppose that $S_{1} \subseteq \mathbb{R}^{n}$ and $S_{2} \subseteq \mathbb{R}^{m}$ are open sets such that

$$
C_{1} \times C_{2} \subset S_{1} \times S_{2}
$$

A differentiable function $f: X \rightarrow R^{k}$ and $G_{f}$ such that every component $G_{f_{i}}$ is strictly increasing on the range of $I_{f_{i}}$.

Definition 6. If $\exists G_{f}$ and $\eta$ such that $\forall \varphi \in X$ and $p_{i} \in R^{n}$, we have

$$
\begin{aligned}
& \left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right. \\
& \left.+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], \ldots, G_{f_{k}}\left(f_{k}(\varphi)\right)-G_{f_{k}}\left(f_{k}(\delta)\right)+\frac{1}{2} p_{k}^{T}\left[G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}\right. \\
& \left.\left.+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right] p_{k}-\eta^{T}(\varphi, \delta)\left[G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)+\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right\} p_{k}\right]\right\} \in K
\end{aligned}
$$

then $f$ is $K-G_{f}$-bonvex at $\delta \in X$ with respect to $\eta$.
Definition 7. If $\exists G_{f}$ and $\eta$ such that $\forall \varphi \in X$ and $p_{i} \in R^{m}$, we have

$$
\begin{aligned}
& \left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right. \\
& \left.+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], \ldots, G_{f_{k}}\left(f_{k}(\varphi)\right)-G_{f_{k}}\left(f_{k}(\delta)\right)+\frac{1}{2} p_{k}^{T}\left[G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}\right. \\
& \left.\left.\quad+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right] p_{k}-\eta^{T}(\varphi, \delta)\left[G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)+\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right\} p_{k}\right]\right\} \in-K,
\end{aligned}
$$

then $f$ is $K-G_{f}$-boncave at $\delta \in X$ with respect to $\eta$.
Generalized the above definitions on two variable, as follows,
Definition 8. If $\exists$ and $G_{f}$ and $\eta_{1}$ such that $\forall \varphi \in X$ and $q_{i} \in R^{n}$, we have

$$
\begin{aligned}
\left\{G_{f_{1}}( \right. & \left.f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} q_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right] q_{1}-\eta_{1}^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right)\right. \\
& \left.\nabla_{\varphi} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\} q_{1}\right], \ldots, G_{f_{k}}\left(f_{k}(\varphi, \ell)\right)-G_{f_{k}}\left(f_{k}(\delta, \ell)\right) \\
& +\frac{1}{2} q_{k}^{T}\left[G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right] q_{k} \\
- & \left.\eta_{1}^{T}(\varphi, \delta)\left[G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta, \ell)+\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right\} q_{k}\right]\right\} \in K,
\end{aligned}
$$

then, $f$ is $K-G_{f}$-bonvex in the first variable at $\delta \in X$ for fixed $\ell \in Y$ with $\eta_{1}$,
and
If $\exists G_{f} \eta_{2}$ such that $\forall \vartheta \in Y$ and $p_{i} \in R^{m}$, we have

$$
\begin{aligned}
&\left\{G_{f_{1}}\left(f_{1}(\delta, \vartheta)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime} f_{1}(\delta, \ell) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell)\right] p_{1}-\eta_{2}^{T}(\ell, \vartheta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right)\right.\right. \\
&\left.\nabla_{\vartheta} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell)\right\} p_{1}\right], \ldots, G_{f_{k}}\left(f_{k}(\delta, \vartheta)\right)-G_{f_{k}}\left(f_{k}(\delta, \ell)\right) \\
&+\frac{1}{2} p_{k}^{T}\left[G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{k}(\delta, \ell)\right] p_{k} \\
&-\left.\eta_{2}^{T}(\ell, \vartheta)\left[G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)+\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{k}(\delta, \ell)\right\} p_{k}\right]\right\} \in K,
\end{aligned}
$$

then, $f$ is $K-G_{f}$-bonvex in the second variable at $\ell \in Y$ for fixed $\delta \in X$ with $\eta_{2}$.

Definition 9. If $\exists G_{f}$ and $\eta_{1}$ such that $\forall \varphi \in X$ and $q_{i} \in R^{n}$, we have

$$
\begin{aligned}
&\left\{G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} q_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right] q_{1}-\eta_{1}^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right)\right.\right. \\
&\left.\nabla_{\varphi} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\} q_{1}\right], \ldots, G_{f_{k}}\left(f_{k}(\varphi, \ell)\right)-G_{f_{k}}\left(f_{k}(\delta, \ell)\right) \\
&+\frac{1}{2} q_{k}^{T}\left[G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right] q_{k} \\
&-\left.\eta_{1}^{T}(\varphi, \delta)\left[G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta, \ell)+\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right\} q_{k}\right]\right\} \in-K,
\end{aligned}
$$

then, $f$ is $K-G_{f}$-boncave in the first variable at $\delta \in X$ for fixed $\ell \in Y$ with respect to $\eta_{1}$,
and
If $\exists G_{f}$ and $\eta_{2}$ such that $\forall \vartheta \in Y$ and $p_{i} \in R^{m}$, we have

$$
\begin{aligned}
&\left\{G_{f_{1}}\right.\left(f_{1}(\delta, \vartheta)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell)\right] p_{1}-\eta_{2}^{T}(\ell, \vartheta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right)\right. \\
&\left.\nabla_{\vartheta} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell)\right\} p_{1}\right], \ldots, G_{f_{k}}\left(f_{k}(\delta, \vartheta)\right)-G_{f_{k}}\left(f_{k}(\delta, \ell)\right) \\
& \quad+\frac{1}{2} p_{k}^{T}\left[G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}\right. \\
&-\left.\eta_{2}^{T}(\ell, \vartheta)\left[G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)+\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{k}(\delta, \ell)\right\} p_{k}\right]\right\} \in-K,
\end{aligned}
$$

then function $f$ is $K-G_{f}$-boncave in the second variable at $\ell \in Y$ for fixed $\delta \in X$ with respect to $\eta_{2}$.
Example 1. Let $X=[1,2] \subseteq \mathbb{R}, n=m=1$ and $k=2$. Consider $f: X \rightarrow \mathbb{R}^{2}$ be defined by

$$
f(\varphi)=\left(f_{1}(\varphi), f_{2}(\varphi)\right)
$$

where,

$$
f_{1}(\varphi)=\varphi \sin \left(\frac{1}{\varphi}\right), f_{2}(\varphi)=\cos \varphi .
$$

Next, $G_{f}:\left(G_{f_{1}}, G_{f_{2}}\right): \mathbb{R} \rightarrow \mathbb{R}^{2}$ defined by

$$
G_{f_{1}}=t^{2}, \quad G_{f_{2}}=t^{4}
$$

Let $K=\{(\varphi, \vartheta) ; \varphi \geqq 0$ and $\vartheta \geqq 0\}$ and $\eta: X \times X \rightarrow \mathbb{R}$ be given by

$$
\eta(\varphi, \delta)=\left(1-\delta^{2}\right) .
$$

Now, we have to claim that $f$ is $K-G_{f}$-bonvex, for this, we have driven that the following expression as

$$
\begin{aligned}
& \left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right. \\
& \left.+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right. \\
& \left.\left.\quad+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right] p_{2}-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} p_{2}\right]\right\} \in K .
\end{aligned}
$$

## Let

$$
\begin{aligned}
\Pi= & \left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right. \\
+\{ & \left.\left.G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right. \\
& \left.\left.+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right] p_{2}-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} p_{2}\right]\right\}
\end{aligned}
$$

Substituting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta$, we obtain

$$
\begin{aligned}
\Pi= & \left\{\varphi^{2} \sin ^{2} \frac{1}{\varphi^{2}}-\delta^{2} \sin ^{2} \frac{1}{\delta^{2}}+\frac{1}{2} p^{2}\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right]-\left(1-\delta^{2}\right)\left[2 \delta \sin \frac{1}{\delta}\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)\right.\right. \\
& \left.+p\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right]\right], \cos ^{4} \varphi-\cos ^{4} \delta+\frac{1}{2} p^{2}\left[12 \cos ^{2} \delta(-\sin \delta)^{2}\right. \\
& \left.\left.+4 \cos ^{3} \delta(-\cos \delta)\right]-\left(1-\delta^{2}\right)\left[4 \cos ^{3} \delta(-\sin \delta)+p\left(12 \cos ^{2} \delta(-\sin \delta)^{2}+4 \cos ^{3} \delta(-\cos \delta)\right)\right]\right\} .
\end{aligned}
$$

Now, we consider

$$
\begin{gathered}
\Psi=\varphi^{2} \sin ^{2} \frac{1}{\varphi^{2}}-\delta^{2} \sin ^{2} \frac{1}{\delta^{2}}+\frac{1}{2} p^{2}\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right] \\
-\left(1-\delta^{2}\right)\left[2 \delta \sin \frac{1}{\delta}\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)+p\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right]\right]
\end{gathered}
$$

Let us apply the following ansatz:

$$
\Psi=\Psi_{1}+\Psi_{2}(s a y)
$$

consider

$$
\begin{gathered}
\Phi=\left\{\cos ^{4} \varphi-\cos ^{4} \delta+\frac{1}{2} p^{2}\left[12 \cos ^{2} \delta(-\sin \delta)^{2}+4 \cos ^{3} \delta(-\cos \delta)\right]\right. \\
\left.-\left(1-\delta^{2}\right)\left[4 \cos ^{3} \delta(-\sin \delta)+p\left(12 \cos ^{2} \delta(-\sin \delta)^{2}+4 \cos ^{3} \delta(-\cos \delta)\right)\right]\right\} \in K .
\end{gathered}
$$

The above expression breaks in $\Phi_{1}$ and $\Phi_{2}$ (say) as follows:

$$
\Phi=\Phi_{1}+\Phi_{2}
$$

where

$$
\Psi_{1}=\varphi^{2} \sin ^{2} \frac{1}{\varphi^{2}}-\delta^{2} \sin ^{2} \frac{1}{\delta^{2}}-\left(1-\delta^{2}\right)\left[2 \delta \sin \frac{1}{\delta}\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)\right]
$$

It is easily verified from Figure 1, we have

$$
\Psi_{1} \geqq 0, \forall \varphi, \delta \in X
$$

$$
\Psi_{2}=\frac{1}{2} p^{2}\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right]+p\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right]
$$



Figure 1. $\Psi_{1}=\left\{\varphi^{2} \sin ^{2} \frac{1}{\varphi^{2}}-\delta^{2} \sin ^{2} \frac{1}{\delta^{2}}-\left(1-\delta^{2}\right)\left[2 \delta \sin \frac{1}{\delta}\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)\right]\right\}$.
It is clear from Figure 2, we obtain

$$
\Psi_{2} \geqq 0, \forall \delta \in X \text { and } p \in\left[-\frac{1}{10^{10}},-1\right]
$$



Figure 2. $\Psi_{2}=\frac{1}{2} p^{2}\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right]+$ $p\left[2\left(\sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}\right)^{2}+2 \delta \sin \frac{1}{\delta}\left(-\frac{1}{\delta^{3}} \sin \frac{1}{\delta}\right)\right]$.

Now,

$$
\Phi_{1}=\cos ^{4} \varphi-\cos ^{4} \delta+-\left(1-\delta^{2}\right)\left[4 \cos ^{3} \delta(-\sin \delta)\right]
$$

as can be seen from Figure 3.

$$
\Phi_{1} \geqq 0 \forall \varphi, \delta \in X
$$



Figure 3. $\Phi_{1}=\left\{\cos ^{4} \varphi-\cos ^{4} \delta+-\left(1-\delta^{2}\right) 4 \cos ^{3} \delta(-\sin \delta)\right\}$.
and

$$
\Phi_{2}=\frac{1}{2} p^{2}\left[12 \cos ^{2} \delta(-\sin \delta)^{2}+4 \cos ^{3} \delta(-\cos \delta)+p\left(12 \cos ^{2} \delta(-\sin \delta)^{2}+4 \cos ^{3} \delta(-\cos \delta)\right)\right]
$$

As can be seen from Figure 4. $\Phi_{2} \geqq 0, \forall \delta \in X$ and $p_{1}, p_{2} \in\left[-\frac{1}{10^{10}},-1\right]$. (From Figure 4).


Figure 4. $\Phi_{2}=\frac{1}{2} p^{2}\left[12 \cos ^{2} \delta(-\sin \delta)^{2}+4 \cos ^{3} \delta(-\cos \delta)+p\left(12 \cos ^{2} \delta(-\sin \delta)^{2}+4 \cos ^{3} \delta(-\cos \delta)\right)\right]$.
Hence, $\Psi \geqq 0$ and $\Phi \geqq 0$. This gives $\psi+\phi \geqq 0$. Thus, we can find that $(\Psi, \Phi) \in K$.
Hence, $f$ is $K-G_{f}$-bonvex function at $(\Psi, \Phi)$ w.r.t. $\eta$.
We will show that $f$ is not invex. For this it is either

$$
f_{1}(\varphi)-f_{1}(\delta)-\eta^{T}(\varphi, \delta) \nabla_{\varphi} f_{1}(\delta) \nsupseteq 0
$$

or

$$
f_{2}(\varphi)-f_{2}(\delta)-\eta^{T}(\varphi, \delta) \nabla_{\varphi} f_{2}(\delta) \nsucceq 0 .
$$

Since $f_{1}(\varphi)-f_{1}(\delta)-\eta^{T}(\varphi, \delta) \nabla_{\varphi} f_{1}(\delta)=\varphi \sin \frac{1}{\varphi}-\delta \sin \frac{1}{\delta}-\left(1-\delta^{2}\right) \sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta} \nsupseteq 0$, is not $\forall \varphi, \delta \in X$ as can be seen from Figure 5. Also, $f_{2}(\varphi)-f_{1}(\delta)-\eta^{T}(\varphi, \delta) \nabla_{\varphi} f_{2}(\delta)=$ $\cos \varphi-\cos \delta+\left(1-\delta^{2}\right) \sin \delta \nsupseteq 0$, is not $\forall \varphi, \delta \in X$ as can be seen from Figure 6 .


Figure 5. $\varphi \sin \frac{1}{\varphi}-\delta \sin \frac{1}{\delta}-\left(1-\delta^{2}\right) \sin \frac{1}{\delta}-\frac{1}{\delta} \cos \frac{1}{\delta}$.


Figure 6. $\cos \varphi-\cos \delta+\left(1-\delta^{2}\right) \sin \delta$.
Therefore, from the above example, it shows that $f$ is $K-G_{f}$-bonvex, but it is not invex with respect to same $\eta$.

Definition 10. If $\exists G_{f}$ and $\eta$ such that $\forall \varphi \in X$ and $q_{i} \in R^{n}$, we have

$$
\begin{aligned}
& \eta^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)+q_{1}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\}, \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)+q_{k}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}\right.\right. \\
& \left.\left.\quad+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right\}\right\} \in K \Rightarrow\left[G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} q_{1}\right. \\
& \left.\quad, \ldots, G_{f_{k}}\left(f_{k}(\varphi)\right)-G_{f_{k}}\left(f_{k}(\delta)\right)+\frac{1}{2} q_{k}^{T}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right\} q_{k}\right] \in K
\end{aligned}
$$

then, $f$ is $G_{f}$-pseudobonvex at $\delta \in X$ with $\eta$.

Definition 11. If $\exists G_{f}$ and $\eta$ such that $\forall \varphi \in X$ and $q_{1} \in R^{n}$, we have

$$
\begin{aligned}
& \eta^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)+q_{1}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\}, \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)+q_{k}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}\right.\right. \\
& \left.\left.\quad+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right\}\right\} \in-K \Rightarrow\left[G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} q_{1}\right. \\
& \left.\quad, \ldots, G_{f_{k}}\left(f_{k}(\varphi)\right)-G_{f_{k}}\left(f_{k}(\delta)\right)+\frac{1}{2} q_{k}^{T}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta)\right) \nabla_{\varphi} f_{k}(\delta)\left(\nabla_{\varphi} f_{k}(\delta)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta)\right) \nabla_{\varphi \varphi} f_{k}(\delta)\right\} q_{k}\right] \in-K,
\end{aligned}
$$

then $f$ is $G_{f}$-pseudoboncave at $\delta \in X$ with respect to $\eta$.
We generalized the above definition as follows:
Definition 12. If $\exists G_{f}$ and $\eta_{1}$ such that $\forall \varphi \in X$ and $q_{i} \in R^{n}$, we have

$$
\begin{aligned}
& \eta_{1}^{T}(\varphi, \delta)\{ \left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)+q_{1}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\}, \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell)\right. \\
&\left.+q_{k}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right\}\right\} \in K \\
& \Rightarrow\left[G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\} q_{1}, \ldots, G_{f_{k}}\left(f_{k}(\varphi, \ell)\right)\right. \\
&\left.\quad-G_{f_{k}}\left(f_{k}(\delta, \ell)\right)+\frac{1}{2} q_{k}^{T}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right\} q_{k}\right] \in K,
\end{aligned}
$$

then $f$ is $K$ - $G_{f}$-bonvex in the first variable at $\delta \in X$ for fixed $\ell \in Y$ with $\eta_{1}$,
and
if $\exists G_{f}$ and $\eta_{2}$ such that $\forall \vartheta \in Y$ and $p_{i} \in R^{m}$, we have

$$
\begin{aligned}
& \eta_{2}^{T}(\delta, \vartheta)\{ G_{f_{1}}^{\prime}\left(f_{1}(\delta, \vartheta)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell) \ell\right\} p_{1}, \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell) \\
&\left.\quad+p_{k}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{k}(\delta, \ell)\right\}\right\} \in K \\
& \Rightarrow\left[G_{f_{1}}\left(f_{1}(\delta, \vartheta)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} p_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell)\right\} p_{1}, \ldots, G_{f_{k}}\left(f_{k}(\delta, \vartheta)\right)\right. \\
&\left.\quad-G_{f_{k}}\left(f_{k}(\delta, \ell)\right)+\frac{1}{2} p_{k}^{T}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{k}(\delta, \ell)\right\} p_{k}\right] \in K,
\end{aligned}
$$

then $f$ is $K-G_{f}$-bonvex in the second variable at $\ell \in Y$ for fixed $\delta \in X$ with $\eta_{2}$.
Definition 13. If $\exists G_{f}$ and $\eta_{1}$ such that $\forall \varphi \in X$ and $q_{i} \in R^{n}$, we have

$$
\begin{aligned}
& \eta_{1}^{T}(\varphi, \delta)\{ G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)+q_{1}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\}, \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell) \\
&\left.+q_{k}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right\}\right\} \in-K \\
& \Rightarrow\left[G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\} q_{1}, \ldots, G_{f_{k}}\left(f_{k}(\varphi, \ell)\right)\right. \\
&\left.\quad-G_{f_{k}}\left(f_{k}(\delta, \ell)\right)+\frac{1}{2} q_{k}^{T}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi} f_{k}(\delta, \ell)\left(\nabla_{\varphi} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{k}(\delta, \ell)\right\} q_{k}\right] \in-K,
\end{aligned}
$$

then $f$ is $K$ - $G_{f}$-bonvex in the first variable at $\delta \in X$ for fixed $\ell \in Y$ with $\eta_{1}$,
and

If $\exists G_{f}$ and $\eta_{2}$ such that $\forall \vartheta \in Y$ and $p_{i} \in R^{m}$, we have

$$
\begin{aligned}
\eta_{2}^{T}(\delta, \vartheta)\{ & G_{f_{1}}^{\prime}\left(f_{1}(\delta, \vartheta)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell)\right\} p_{1}, \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell) \\
& \left.+p_{k}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{k}(\delta, \ell)\right\}\right\} \in-K \\
\Rightarrow & {\left[G_{f_{1}}\left(f_{1}(\delta, \vartheta)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} p_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta} f_{1}(\delta, \ell)\left(\nabla_{\vartheta} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{1}(\delta, \ell)\right\} p_{1}, \ldots, G_{f_{k}}\left(f_{k}(\delta, \vartheta)\right)\right.} \\
& \left.\quad-G_{f_{k}}\left(f_{k}(\delta, \ell)\right)+\frac{1}{2} p_{k}^{T}\left\{G_{f_{k}}^{\prime \prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta} f_{k}(\delta, \ell)\left(\nabla_{\vartheta} f_{k}(\delta, \ell)\right)^{T}+G_{f_{k}}^{\prime}\left(f_{k}(\delta, \ell)\right) \nabla_{\vartheta \vartheta} f_{k}(\delta, \ell)\right\} p_{k}\right] \in-K .
\end{aligned}
$$

then $f$ is $K-G_{f}$-boncave in the second variable at $\ell \in Y$ for fixed $\delta \in X$ with respect to $\eta_{2}$.

Remark 1. If $G_{f}(t)=t$, then above definition reduces in $K-\eta$-pseudo bonvex w.r.t. $\eta$,

$$
\begin{aligned}
& \eta^{T}(\varphi, \delta)\left[\nabla_{\varphi} f_{1}(\delta)+\nabla_{\varphi \varphi} f_{1}(\delta) q_{1}, \ldots . ., \nabla_{\varphi} f_{k}(\delta)+\nabla_{\varphi \varphi} f_{k}(\delta) q_{k}\right] \in K \\
& \quad \Rightarrow\left[f_{1}(\varphi)-f_{1}(\delta)+\frac{1}{2} q_{1}^{T} \nabla_{\varphi \varphi} f_{1}(\delta) q_{1}, \ldots ., f_{k}(\varphi)-f_{k}(\delta)+\frac{1}{2} q^{T} \nabla_{\varphi \varphi} f_{k}(\delta) q_{k}\right] \in K .
\end{aligned}
$$

Example 2. Let $X=[-10,10]$ and $K=\{(\varphi, \vartheta): \varphi \geqq 0, \varphi \leqq \vartheta\}$. Consider the function $f: X \rightarrow \mathbb{R}^{2}$ defined by

$$
f(\varphi)=\left(f_{1}(\varphi), f_{2}(\varphi)\right)
$$

where

$$
f_{1}(\varphi)=\sin \varphi, \quad f_{2}(\varphi)=e^{\varphi}
$$

Define $G_{f}=\left(G_{f_{1}}, G_{f_{2}}\right): R^{2} \rightarrow R$ given by

$$
G_{f_{1}}=t^{2}, G_{f_{2}}=t^{3}, \eta=\varphi^{2}-\delta^{2}, \text { and } q_{1}=q_{2} \in[2, \infty] .
$$

We have to claim that function $f$ is $K-G_{f}$-pseudobonvex at point $\delta$, i.e.,

$$
\begin{gathered}
\eta^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)+q_{1}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\}, G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)+q_{2}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right.\right. \\
\left.\left.+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\}\right\} \in K \Rightarrow\left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} q_{1},\right. \\
\left.G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} q_{2}^{T}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} q_{2}\right\} \in K . \\
\text { Consider } \\
\quad \\
\left.\quad+q_{2}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\}\right\} .
\end{gathered}
$$

Putting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta$, we have

$$
\tau=\left(\varphi^{2}-\delta^{2}\right)\left(\sin 2 \delta+2 q_{1}\left(\cos \delta-\sin ^{2} \delta\right), 3 e^{3 \delta}+9 e^{2 \delta} q_{2}\right)
$$

At the point $\delta=0$, the value of above expression becomes

$$
\tau=\left\{2 \varphi^{2} q_{1}, 3 \varphi^{2}\left(1+3 q_{2}\right)\right\}, \quad \forall q_{1}=q_{2} \in[2, \infty)
$$

Obviously,

$$
\tau=\left\{2 \varphi^{2} q_{1}, 3 \varphi^{2}\left(1+3 q_{2}\right)\right\} \in K .
$$

Next, consider

$$
\begin{gathered}
\Psi=\left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} q_{1}, \quad G_{f_{2}}\left(f_{2}(\varphi)\right)\right. \\
\left.-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} q_{2}^{T}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} q_{2}\right\} .
\end{gathered}
$$

Putting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta$, we have

$$
\Psi=\left\{\sin ^{2} \varphi-\sin ^{2} \delta+\frac{1}{2} q_{1}^{2}\left(2 \cos ^{2} \delta-2 \sin ^{2} \delta\right), e^{3 \varphi}-e^{3 \delta}+\frac{9}{2} q_{2}^{2} e^{3 \delta}\right\} .
$$

The value of above expression at the point $\delta=0$, we get

$$
\Psi=\left\{\sin ^{2} \varphi+q_{1}^{2}, e^{3 \varphi}+\frac{9}{2} q_{2}^{2}-1\right\} \in K .
$$

From the Figure 7. We can easily observe that the value of $\varphi$-coordinate always less than $\vartheta$-coordinate in $K$, so $\varphi \in K$.

Hence, $f$ is $K-G_{f}$-pseudobonvex at the point $\delta=0$ with respect to $\eta$.
Next,

$$
\begin{aligned}
& \left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right. \\
& \left.+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], \\
& G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right. \\
& \left.\left.+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right] p_{2}-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi} f_{2}(\delta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} p_{2}\right]\right\} \notin K .
\end{aligned}
$$

Let

$$
\begin{aligned}
\Psi= & \left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right. \\
+ & \left.\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right. \\
& \left.\left.+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right] p_{2}-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi} f_{2}(\delta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} p_{2}\right]\right\} .
\end{aligned}
$$

Substituting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta$, we obtain
$\Psi=\left\{\sin ^{2} \varphi-\sin ^{2} \delta+p_{1}^{2}\left(\cos ^{2} \delta-\sin ^{2} \delta\right)-\left(\varphi^{2}-\delta^{2}\right)\left(\sin 2 \delta+2 p_{1}\left(\cos \delta-\sin ^{2} \delta\right)\right), e^{3 \varphi}-e^{3 \delta}+\frac{9}{2} p_{2}^{2} e^{3 \delta}-\left(\varphi^{2}-\delta^{2}\right)\left(3 e^{3 \delta}+9 e^{2 \delta} p_{2}\right)\right\}$.
At the point $\delta=0$, it follows that

$$
\Psi=\left\{\sin ^{2} \varphi+p_{1}^{2}-2 p_{1} \varphi^{2}, e^{3 \varphi}+\frac{9}{2} p_{2}^{2}-1-\varphi^{2}\left(3+9 p_{2}\right)\right\}, p_{1}=p_{2} \in[2, \infty)
$$

Take particular point $\varphi=-\frac{\pi}{2}$ and $p_{1}=p_{2}=2 \in[2, \infty)$, we obtain,

$$
\Psi=(-4.86,-34.80) \notin K .
$$

Hence, $f$ is $K-G_{f}$-pseudobonvex, but it is not $K-G_{f}$-bonvex at $\delta=0$ with respect to $\eta$.


Figure 7. $\left(\sin ^{2} \varphi+4 \cos ^{2} \varphi, e^{3 \varphi}+17\right)$.

In the following example, we showed that the function $f$ is $K-G_{f}$-pseudobonvex, but it is not $K-G_{f}$-bonvex function with same $\eta$.

Example 3. Let $X=\left[0, \frac{\pi}{2}\right]$ and $K=\{(\varphi, \vartheta): \varphi \geqq 0, \vartheta \geqq \varphi\}$. Consider $G_{f}=\left(G_{f_{1}}, G_{f_{2}}\right)$ : $R^{2} \rightarrow R$ and $f: X \rightarrow \mathbb{R}^{2}$ given by

$$
f(\varphi)=\left(f_{1}(\varphi), f_{2}(\varphi)\right)
$$

where

$$
\begin{gathered}
f_{1}(\varphi)=\sin \varphi, \quad f_{2}(\varphi)=\varphi, \\
G_{f_{1}}=t, \quad G_{f_{2}}=t^{2}
\end{gathered}
$$

Define $\eta: X \times X \rightarrow R^{n}$ given by

$$
\eta(\varphi, \delta)=\varphi-\delta \text { and } q_{1}, q_{2} \in[1, \infty] .
$$

Solution: In this example, we will try to derive that $f$ is $K-G_{f}$-pseudobonvex i.e.,

$$
\begin{aligned}
& \eta^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)+q_{1}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\}, G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)+q_{2}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right.\right. \\
& \left.\left.+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\}\right\} \in K \Rightarrow\left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} q_{1},\right. \\
& \left.\quad G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} q_{2}^{T}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} q_{2}\right\} \in K .
\end{aligned}
$$

Consider

$$
\begin{aligned}
\Pi_{1}=\eta^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right. & +q_{1}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\}, G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta) \\
& \left.+q_{2}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\}\right\} .
\end{aligned}
$$

Putting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta$, we have

$$
\Pi_{1}=\left\{(\varphi-\delta) \cos \delta,(\varphi-\delta)\left(2 \delta+2 q_{2}\right)\right\}
$$

The value of above expression at the point $\delta=0$, we get

$$
\Pi_{1}=\left\{\varphi, 2 \delta q_{2}\right\} \in K
$$

Next, let

$$
\begin{gathered}
\Pi_{2}=\left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} q_{1}^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} q_{1}, \quad G_{f_{2}}\left(f_{2}(\varphi)\right)\right. \\
\left.-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} q_{2}^{T}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta)\right) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} q_{2}\right\} .
\end{gathered}
$$

Putting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta$, we have

$$
\Pi_{2}=\left\{\sin \varphi-\sin \delta+\frac{1}{2} q_{1}^{2}(-\sin \delta), \varphi-\delta+q_{2}^{2}\right\} .
$$

After simplifying and the value at $\delta=0$, it follows that

$$
\Pi_{2}=\left\{\sin \varphi, \varphi+q_{2}^{2}\right\} \in K .
$$

Hence, $f$ is $K-G_{f}$-pseudobonvex at the point $\delta=0$ with respect to $\eta$.

Next,

$$
\begin{aligned}
& \left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right. \\
& \left.+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], \\
& G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right. \\
& \left.\left.+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right] p_{2}-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi} f_{2}(\delta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} p_{2}\right]\right\} \notin K .
\end{aligned}
$$

Let
$\Psi=\left\{G_{f_{1}}\left(f_{1}(\varphi)\right)-G_{f_{1}}\left(f_{1}(\delta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right] p_{1}-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\right.\right.$ $\left.+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta)\right) \nabla_{\varphi} f_{1}(\delta)\left(\nabla_{\varphi} f_{1}(\delta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta)\right) \nabla_{\varphi \varphi} f_{1}(\delta)\right\} p_{1}\right], G_{f_{2}}\left(f_{2}(\varphi)\right)-G_{f_{2}}\left(f_{2}(\delta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}\right.$ $\left.\left.+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right] p_{2}-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi} f_{2}(\delta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta)\right) \nabla_{\varphi} f_{2}(\delta)\left(\nabla_{\varphi} f_{2}(\delta)\right)^{T}+G_{f_{2}}^{\prime} f_{2}(\delta) \nabla_{\varphi \varphi} f_{2}(\delta)\right\} p_{2}\right]\right\}$.

Substituting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta$, we obtain

$$
\Psi=\left\{\sin \varphi-\sin \delta+\frac{1}{2} p_{1}^{2}(-\sin \delta)-(\varphi-\delta) p_{1} \cos \delta, \varphi^{2}+p_{2}^{2}-2(\varphi-\delta) p_{2}\right\}
$$

At the point $\delta=0$, it follows that

$$
\Psi=\left\{\sin \varphi-p_{1} \varphi,\left(\varphi-p_{2}\right)^{2}\right\} \notin K .
$$

Hence, $f$ is $K-G_{f}$-pseudobonvex, but it is not $K-G_{f}$-bonvex at $\delta=0$ with respect to $\eta$.

## 3. $K-G_{f}$-Wolfe Type Second-Order Symmetric Primal-Dual Pair with Cones

The study of second-order duality is more significant due to computational advantage over first order duality as it provides tighter bounds for the objective functions, when approximation is used.

The motivated by [21-27] several researches in this area, we formulated a new type $K-G_{f}$-Wolfe type primal dual pair, with cone objectives as well as cone constraint as follows: Primal Problem (GWPP):
$K$-min $L(\varphi, \vartheta, \lambda, p)=\left\{L_{1}(\varphi, \vartheta, \lambda, p), L_{2}(\varphi, \vartheta, \lambda, p), L_{3}(\varphi, \vartheta, \lambda, p), \ldots, L_{k}(\varphi, \vartheta, \lambda, p)\right\}$, where

$$
\begin{aligned}
& L_{i}(\varphi, \vartheta, \lambda, p)=G_{f_{i}}\left(f_{i}(\varphi, \vartheta)\right)-\vartheta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}\right.\right. \\
& \left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\varphi, \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} p_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i},
\end{aligned}
$$

subject to

$$
\begin{gather*}
-\sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}\right] \in C_{2}^{*}  \tag{1}\\
\lambda^{T} e_{k}=1, \lambda \in \operatorname{int} K^{*}, \quad \varphi \in C_{1} \tag{2}
\end{gather*}
$$

## Dual Problem (GWDP):

$K-\max M(\delta, \ell, \lambda, q)=\left\{M_{1}(\delta, \ell, q), M_{2}(\delta, \ell, \lambda, q), M_{3}(\delta, \ell, \lambda, q), \ldots, M_{k}(\delta, \ell, \lambda, q)\right\}$,
where

$$
\begin{aligned}
& M_{i}(\delta, \ell, \lambda, q)=G_{f_{i}}\left(f_{i}(\delta, \ell)\right)-\delta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}\right.\right. \\
& \left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} q_{i}\left[G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right] q_{i},
\end{aligned}
$$

subject to

$$
\begin{equation*}
\sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right] \in C_{1}^{*} \tag{3}
\end{equation*}
$$

$$
\begin{equation*}
\lambda^{T} e_{k}=1, \quad \lambda \in \operatorname{int} K^{*}, \quad \delta \in C_{2} \tag{4}
\end{equation*}
$$

where, for $i \in \tilde{Q}$,

- $f_{i}: R_{1} \times R_{2} \rightarrow R$, is a differential function of $\varphi$ and $\vartheta, e_{k}=(1,1, \ldots, 1)^{T} \in R^{k}$,
- $\quad q_{i}$ and $p_{i}$ are vectors in $R^{n}$ and $R^{m}$, respectively and $\lambda \in R^{k}$.

Let $V^{*}$ and $W^{*}$ be the sets of feasible solutions of (GWPP) and (GWDP) respectively.
Theorem 1 (Weak duality). Let $(\varphi, \vartheta, \lambda, p) \in V^{*}$ and $(\delta, \ell, \lambda, q) \in W^{*}$. Let, for $i \in \tilde{N}$
(i) $\left\{f_{1}(., \ell), f_{2}(., \ell), \ldots, f_{k}(., \ell)\right\}$ be $K-G_{f_{i}}$-bonvex at $\delta$ w.r.t. $\eta_{1}$,
(ii) $\left\{f_{1}(\varphi,),. f_{2}(\varphi,),. \ldots, f_{k}(\varphi,).\right\}$ be $K-G_{f_{i}}$-boncave in $\vartheta$ w.r.t. $\eta_{2}$,
(iii) $\quad \eta_{1}(\varphi, \delta)+\delta \in C_{1}, \forall(\varphi, \delta) \in C_{1} \times C_{2}$,
(iv) $\eta_{2}(\ell, \vartheta)+\vartheta \in C_{2}, \forall(\ell, \vartheta) \in C_{1} \times C_{2}$,

Then, $L(\varphi, \vartheta, \lambda, p)-M(\delta, \ell, \lambda, q) \notin-K \backslash\{0\}$.
Proof. If possible, then suppose

$$
L(\varphi, \vartheta, \lambda, p)-M(\delta, \ell, \lambda, q) \in-K \backslash\{0\}
$$

or
$\left\{G_{f_{1}}\left(f_{1}(\varphi, \vartheta)\right)-\vartheta^{T} \sum_{i=1}^{k} \lambda_{i}\left(G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}\right)\right.$
$-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} p_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}, \ldots, G_{f_{k}}\left(f_{k}(\varphi, \vartheta)\right)$
$-\vartheta^{T} \sum_{i=1}^{k} \lambda_{i}\left(G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}\right)$
$-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} p_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)-\delta^{T} \sum_{i=1}^{k} \lambda_{i}\left(G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right)\right.$
$\left.\nabla_{\varphi} f_{i}(\delta, \ell)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right)-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} q_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}\right.$
$\left.+G_{f_{i}}^{\prime}(f(\delta, \ell)) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}, \ldots, G_{f_{k}}\left(f_{k}(\delta, \ell)\right)-\delta^{T} \sum_{i=1}^{k} \lambda_{i}\left(G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}\right.\right.$
$\left.\left.\left.+G_{f_{i}^{\prime}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right)-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} q_{i}^{T}\left\{G_{f_{i}^{\prime}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}(f(\delta, \ell)) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right\} \in-K \backslash\{0\}$.
Since $\lambda \in \operatorname{int} K^{*}$, we get

$$
\begin{align*}
& \sum_{i=1}^{k} \lambda_{i}\left\{G_{f_{i}}\left(f_{i}(\varphi, \vartheta)\right)-\vartheta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}\right.\right.\right. \\
& \left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}\right]-\delta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)\right.\right.  \tag{5}\\
& \left.\left.\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}\right\}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} p_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{i}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\}-\left\{G_{f_{i}}\left(f_{i}(\delta, \ell)\right)\right. \\
& \left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\}-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} q_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\varphi, \vartheta)\right\} q_{i}\right\}<0 .
\end{align*}
$$

By hypothesis $(i)$ and using $\lambda \in \operatorname{int} K^{*}$, we get

$$
\begin{aligned}
& \sum_{i=1}^{k} \lambda_{i}\left\{G_{f_{i}}\left(f_{i}(\varphi, \ell)\right)-G_{f_{i}}\left(f_{i}(\delta, \ell)\right)+\frac{1}{2} q_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right. \\
& \left.-\eta_{1}^{T}(\varphi, \delta)\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right]\right\} \geqq 0
\end{aligned}
$$

Using feasibility of dual problem (GWDP) \& using dual constraints with assumption (iii), it yields

$$
\left(\eta_{1}(\varphi, \delta)+\delta\right)^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi} f_{i}(\delta, \ell)\right\} q_{i}\right] \geqq 0
$$

it implies that

$$
\begin{align*}
& \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}\left(f_{i}(\varphi, \ell)\right)-G_{f_{i}}\left(f_{i}(\delta, \ell)\right)+\frac{1}{2} q_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right] \\
& \geqq-\delta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi}\left(f_{i}(\delta, \ell)\right)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} q_{i}\right] \tag{6}
\end{align*}
$$

Similarly, using hypotheses (ii), (iv), feasible conditions of primal problem (GWPP), dual constraint and $\lambda \in \operatorname{int} K^{*}$,
we get

$$
\begin{align*}
& \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}\left(f_{i}(\varphi, \vartheta)\right)-G_{f_{i}}\left(f_{i}(\varphi, \ell)\right)+\frac{1}{2} p_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\vartheta} f_{i}(\delta, \ell)\right)\left(\nabla_{\vartheta} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} p_{i}\right] \\
& \quad \geqq \vartheta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\vartheta}\left(f_{i}(\delta, \ell)\right)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\} p_{i}\right] \tag{7}
\end{align*}
$$

Now, from inequalities (6), (7) and using the fact that $\lambda^{T} e_{k}=1$, we find that

$$
\begin{aligned}
& \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}\left(f_{i}(\varphi, \vartheta)\right)-\vartheta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta}\left(f_{i}(\varphi, \vartheta)\right)+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\delta, \ell)\right)\left(\nabla_{\vartheta} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\} p_{i}\right]\right. \\
& -\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} p_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{i}(\varphi, \vartheta)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{i}(\varphi, \vartheta)\right\}-G_{f_{i}}\left(f_{i}(\delta, \ell)\right) \\
& -\delta^{T} \sum_{i=1}^{k} \lambda_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi}\left(f_{i}(\delta, \ell)\right)+\left\{G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{i}(\delta, \ell)\right\}\right] \\
& \left.-\frac{1}{2} \delta^{T} \sum_{i=1}^{k} \lambda_{i} q_{i}^{T}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{i}(\delta, \ell)\right) q_{i}\right\}\right] \geqq 0
\end{aligned}
$$

we arrive at contradiction.
Through following example, we validate the Weak duality theorem as:
Example 4. Let $n=m=1, k=2, X=[1,2], p \in\left[2^{2}, 2^{10}\right], q \in\left[10^{-19}, 10^{19}\right], K=\{(\varphi, \vartheta) ; \varphi \geqq$ $0, \varphi \geqq \vartheta\}$ and
$-K=\{(\varphi, \vartheta) ; \varphi \leq 0, \varphi \leq \vartheta\}, R_{1}=R_{2}=R_{+}$. Let $f_{i}: R_{1} \times R_{2} \rightarrow R$ and $G_{f_{i}}$ for $i=1,2$. be defined as

$$
f_{1}(\varphi, \vartheta)=\varphi+\cos \vartheta, f_{2}(\varphi, \vartheta)=\sin \vartheta, G_{f_{1}}(t)=t^{2}, G_{f_{2}}(t)=t
$$

Further, let

$$
\eta_{1}(\varphi, \delta)=\varphi \delta, \quad \eta_{2}(\ell, \vartheta)=\ell-\vartheta
$$

Assume that $C_{1}=C_{2}=C_{1}^{*}=C_{2}^{*}=R_{+}$.
(GWPP) K-minimize $L(\varphi, \vartheta, \lambda, p)=\left\{L_{1}(\varphi, \vartheta, \lambda, p), L_{2}(\varphi, \vartheta, \lambda, p)\right\}$

## Subject to constraints

$$
\begin{gather*}
\lambda_{1}\left[2(\varphi+\cos \vartheta)(-\sin \vartheta)+\left\{2 \sin ^{2} \vartheta+2(\varphi+\cos \vartheta)(-\cos \vartheta)\right\} p_{1}\right]+\lambda_{2}\left[\cos \vartheta-p_{2} \sin \vartheta\right] \leqq 0,  \tag{8}\\
\lambda_{1}+\lambda_{2}=1, \lambda_{i} \in \operatorname{int} K^{*}, \varphi \in C_{1}, i=1,2 . \tag{9}
\end{gather*}
$$

(GWDP) K-maximize $M(\delta, \ell, \lambda, q)=\left\{M_{1}(\delta, \ell, \lambda, q), M_{2}(\delta, \ell, \lambda, q)\right\}$
Subject to constraints

$$
\begin{gather*}
\lambda_{1}\left[2(\varphi+\cos \vartheta)+2 q_{1}\right] \geqq 0  \tag{10}\\
\lambda_{1}+\lambda_{2}=1, \quad \lambda_{i} \in \operatorname{int} K^{*}, \varphi \in C_{2}, i=1,2 .
\end{gather*}
$$

(A1). $\left\{f_{1}(., \ell), f_{2}(., \ell)\right\}$ is $K-G_{f}$-bonvex at $\delta=0$ w.r.t. $\eta_{1}, \forall \varphi \in S_{1}$, i.e.,

$$
\left\{G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right] p_{1}\right.
$$

$$
-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\} p_{1}\right],
$$

$$
G_{f_{2}}\left(f_{2}(\varphi, \ell)\right)-G_{f_{2}}\left(f_{2}(\delta, \ell)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)\left(\nabla_{\varphi} f_{2}(\delta, \ell)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{2}(\delta, \ell)\right] p_{2}
$$

$$
\begin{equation*}
\left.-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)\left(\nabla_{\varphi} f_{2}(\delta, \ell)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{2}(\delta, \ell)\right\} p_{2}\right]\right\} \in K . \tag{12}
\end{equation*}
$$

Consider

$$
\begin{align*}
& \Psi=\left\{G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right] p_{1}\right. \\
&-\eta^{T}(\varphi, \delta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\} p_{1}\right] . \\
& G_{f_{2}}\left(f_{2}(\varphi, \ell)\right)-G_{f_{2}}\left(f_{2}(\delta, \ell)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)\left(\nabla_{\varphi} f_{2}(\delta, \ell)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{2}(\delta, \ell)\right] p_{2} \\
&\left.-\eta^{T}(\varphi, \delta)\left[G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)\left(\nabla_{\varphi} f_{2}(\delta, \ell)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{2}(\delta, \ell)\right\} p_{2}\right]\right\} . \tag{1}
\end{align*}
$$

Putting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta_{1}$ at the point $\delta=0$, and simplifying, we get

$$
\Psi=\left(\varphi^{2}+2 \varphi \cos \ell+p^{2}, 0\right) .
$$

It is clear that

$$
\Psi=\left(\varphi^{2}+2 \varphi \cos \ell+p^{2}, 0\right) \in K .
$$

(A2). $\left\{f_{1}(\varphi,),. f_{2}(\varphi,).\right\}$ is $K-G_{f}$-boncave at $\vartheta=0$ w.r.t. $\eta_{2}, \ell \in S_{2}$,
$\left\{G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\varphi, \vartheta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}^{\prime \prime}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right] p_{1}\right.$
$-\eta^{T}(\ell, \vartheta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right\} p_{1}\right]$, $G_{f_{2}}\left(f_{2}(\varphi, \ell)\right)-G_{f_{2}}\left(f_{2}(\varphi, \vartheta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right] p_{2}$

$$
\begin{equation*}
\left.-\eta^{T}(\ell, \vartheta)\left[G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right\} p_{2}\right]\right\} \in-K . \tag{14}
\end{equation*}
$$

Let $\Psi_{1}=\left\{G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\varphi, \vartheta)\right)+\frac{1}{2} p_{1}^{T}\left[G_{f_{1}}^{\prime \prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right] p_{1}\right.$ $-\eta^{T}(\ell, \vartheta)\left[G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)+\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right\} p_{1}\right]$, $G_{f_{2}}\left(f_{2}(\varphi, \ell)\right)-G_{f_{2}}\left(f_{2}(\varphi, \vartheta)\right)+\frac{1}{2} p_{2}^{T}\left[G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right] p_{2}$

$$
\begin{equation*}
-\eta^{T}(\ell, \vartheta)\left[G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)+\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right\} p_{2}\right]\right\} \tag{15}
\end{equation*}
$$

Putting the values of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta_{2}$ at $\vartheta=0$, we obtain

$$
\begin{gathered}
\Psi_{1}=\left((\varphi+\cos \ell)^{2}-(\varphi+1)^{2}-p_{1}^{2}(\varphi+1)+2 \ell(\varphi+1), \sin \ell-\ell\right) . \\
\Psi_{1}=\left((\varphi+\cos \ell)^{2}-(\varphi+1)^{2}-p_{1}^{2}(\varphi+1)+2 \ell(\varphi+1), \sin \ell-\ell\right) \in-K .
\end{gathered}
$$

(A3). $\eta_{1}(\varphi, \delta)+\delta \in C_{1}, \forall \varphi \in C_{1}$.
(A4). $\quad \eta_{2}(\ell, \vartheta)+\vartheta \in C_{2}, \quad \forall \ell \in C_{2}$.
Validation: To validate Weak duality theorem it is enough to claim that any point ( $\varphi, 0, \lambda_{1}, \lambda_{2}, p$ ) such that $\varphi \geqq 0, \lambda_{1}+\lambda_{2}=1$ are feasible to $(G W P P)$. Also, the points $\left(0, \ell, \lambda_{1}, \lambda_{2}, q\right)$ such that $\ell \geqq 0, \lambda_{1}+\lambda_{2}=1$ are feasible to (GWDP). Now, at these feasible points,

$$
L=\left(L_{1}, L_{2}\right)=\left((\varphi+1)^{2}+\lambda_{1} p_{1}^{2}(\varphi+1), \lambda_{1} p_{1}^{2}(\varphi+1)\right)
$$

and

$$
M=\left(M_{1}, M_{2}\right)=\left(\cos ^{2} \ell-\lambda_{1} q_{1}^{2}, \sin \ell-\lambda_{1} q_{1}^{2}\right)
$$

Now, calculate the value at above feasible points, we have

$$
\begin{equation*}
L(\varphi, \vartheta, \lambda, p)-M(\delta, \ell, \lambda, q)=\left((\varphi+1)^{2}+\lambda_{1} p_{1}^{2}(\varphi+1)-\cos ^{2} \ell+\lambda_{1} q_{1}^{2}, \lambda_{1} p_{1}^{2}(\varphi+1)-\sin \ell+\lambda_{1} q_{1}^{2}\right) \notin K \backslash\{0\} . \tag{16}
\end{equation*}
$$

In particular, the points $\left(\varphi, \vartheta, \lambda_{1}, \lambda_{2}, p\right)=\left(1,0, \frac{1}{2}, \frac{1}{2}, 4\right)$ and $\left(\delta, \ell, \lambda_{1}, \lambda_{2}, q\right)$
$=\left(0, \frac{22}{14}, \frac{1}{2}, \frac{1}{2}, 2\right)$ are feasible solutions for (GWPP) and (GWDP), respectively. Also

$$
\begin{equation*}
L(\varphi, \vartheta, \lambda, p)-M(\delta, \ell, \lambda, q)=(22,17) \notin-K \backslash\{0\} \tag{17}
\end{equation*}
$$

Hence, this validate the results.
Remark 2. Every pseudoconvex function is convex function. On the same pattern we can proof that $K-G_{f}$-pseudobonvex is $K-G_{f}$-bonvex with respect to same $\eta$. So, above proof of Weak duality 3.2 follows on same pattern as Theorem 1.

Theorem 2 (Weak duality). Let $(\varphi, \vartheta, \lambda, p) \in V^{*}$ and $(\delta, \ell, \lambda, q) \in W^{*}$. Let, For $i \in \tilde{N}$
(i) $\left\{f_{1}(., \ell), f_{2}(., \ell), \ldots, f_{k}(., \ell)\right\}$ be K-G $G_{f}$-pseudobonvex at $\ell$ w.r.t. $\eta_{1}$,
(ii) $\left\{f_{1}(\varphi,),. f_{2}(\varphi,),. \ldots, f_{k}(\varphi,).\right\}$ be $K-G_{f}$-pseudoboncave at $\vartheta$, w.r.t. $\eta_{2}$,
(iii) $\quad \eta_{1}(\varphi, \delta)+\delta \in C_{1}, \forall(\varphi, \delta) \in C_{1} \times C_{2}$,
(iv) $\eta_{2}(\ell, \vartheta)+\vartheta \in C_{2}, \forall(\ell, \vartheta) \in C_{1} \times C_{2}$,

Then, $L(\varphi, \vartheta, \lambda, p)-M(\delta, \ell, \lambda, q) \notin-K \backslash\{0\}$.
Proof. Proof follows on same lines as Weak Duality Theorem 1.
Example 5. For $n=m=1, k=2, X=[2,3], p \in[0,1], q \in\left[2,2^{10}\right], \quad K=\{(\varphi, \vartheta)$; $\varphi \leqq 0, \vartheta \geqq 0,|\varphi| \geqq \vartheta\}$,

$$
\begin{aligned}
& R_{1}=R_{2}=R_{+} . \text {Let } f_{i}: R_{1} \times R_{2} \rightarrow R \text { be given as } \\
& f_{1}(\varphi, \vartheta)=\varphi+\vartheta^{2}, f_{2}(\varphi, \vartheta)=1-\vartheta, G_{f_{1}}(t)=t^{2}, G_{f_{2}}(t)=t .
\end{aligned}
$$

Further, Let

$$
\eta_{1}(\varphi, \delta)=\varphi \delta, \eta_{2}(\ell, \vartheta)=\ell-\vartheta .
$$

Assume that $C_{1}=C_{2}=C_{1}^{*}=C_{2}^{*}=R_{+}$.
(GWPP) K-minimize $L(\varphi, \vartheta, \lambda, p)=\left\{L_{1}(\varphi, \vartheta, \lambda, p), L_{2}(\varphi, \vartheta, \lambda, p)\right\}$
Subject to constraints

$$
\begin{gather*}
\lambda_{1}\left[4 \vartheta\left(\varphi+\vartheta^{2}\right)+p_{1}\left\{8 \vartheta^{2}+4\left(\varphi+\vartheta^{2}\right)\right\}\right]-\lambda_{2} \leqq 0  \tag{18}\\
\lambda_{1}+\lambda_{2}=1, \quad \lambda_{i} \in \operatorname{int} K^{*}, \varphi \in C_{1}, \quad i=1,2 \tag{19}
\end{gather*}
$$

(GWDP) K-maximize $M(\delta, \ell, \lambda, q)=\left\{M_{1}(\delta, \ell, \lambda, q), M_{2}(\delta, \ell, \lambda, q)\right\}$
Subject to constraints

$$
\begin{gather*}
\lambda_{1}\left[2\left(\delta+\ell^{2}+q\right)\right] \geqq 0  \tag{20}\\
\lambda_{1}+\lambda_{2}=1, \quad \lambda_{i} \in \operatorname{int} K^{*}, \quad \delta \in C_{2}, \quad i=1,2 \tag{21}
\end{gather*}
$$

(A1). $\left\{f_{1}(., \ell), f_{2}(., \ell)\right\}$ is $K-G_{f}$-pseudobonvex at $\delta$ with respect to $\eta_{1}, \varphi \in R_{1}$, so that

$$
\begin{align*}
& \eta_{1}^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)+p\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\}\right. \\
& \left.G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)+p\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{2}(\delta, \ell)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{2}(\delta, \ell)\right\}\right\} \in K \tag{22}
\end{align*}
$$

Let

$$
\begin{gather*}
\Pi_{1}=\eta_{1}^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)+p\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\}\right. \\
\left.G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)+p\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta, \ell)\right)\left(\nabla_{\varphi} f_{2}(\delta, \ell)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{2}(\delta, \ell)\right\}\right\} \tag{23}
\end{gather*}
$$

Next, let

$$
\begin{align*}
\Pi_{2}= & {\left[G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\delta, \ell)\right)+\frac{1}{2} p^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi} f_{1}(\delta, \ell)\left(\nabla_{\varphi} f_{1}(\delta, \ell)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{1}(\delta, \ell)\right\} p\right.} \\
& \left.G_{f_{2}}\left(f_{2}(\varphi, \ell)\right)-G_{f_{2}}\left(f_{2}(\delta, \ell)\right)+\frac{1}{2} p^{T}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi} f_{2}(\delta, \ell)\left(\nabla_{\varphi} f_{2}(\delta, \ell)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\delta, \ell)\right) \nabla_{\varphi \varphi} f_{2}(\delta, \ell)\right\} p\right] . \tag{24}
\end{align*}
$$

After simplification, substituting the value of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta_{1}$ at $\delta=0$, we get

$$
\Pi_{1}=(0,0) \in K \Rightarrow \Pi_{2}=\left(\varphi^{2}-2 \varphi \ell^{2}+p^{2}, 0\right) \in K .
$$

(A2). $\left\{f_{1}(\varphi,),. f_{2}(\varphi,).\right\}$ is $K-G_{f}$-pseudoboncave at $\vartheta$ with respect to $\eta_{2}$ for fixed $\varphi$ for all $\ell \in S_{2}$, i.e.,

$$
\begin{align*}
& \eta_{2}^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)+q\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right\}, G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right. \\
& \left.+q\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right\}\right\} \in K \\
& \Rightarrow\left[G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\varphi, \vartheta)\right)+\frac{1}{2} q^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right\} q,\right. \\
& \left.\quad G_{f_{2}}\left(f_{2}(\varphi, \ell)\right)-G_{f_{2}}\left(f_{2}(\varphi, \vartheta)\right)+\frac{1}{2} q^{T}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right\} q\right] \in-K .  \tag{25}\\
& \text { Let } \quad \Pi_{3}=\eta_{2}^{T}(\varphi, \delta)\left\{G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)+q\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right\},\right. \\
& \left.\quad G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)+q\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right\}\right\}, \\
& \text { and } \\
& \Pi_{4}=\left[G_{f_{1}}\left(f_{1}(\varphi, \ell)\right)-G_{f_{1}}\left(f_{1}(\varphi, \vartheta)\right)+\frac{1}{2} q^{T}\left\{G_{f_{1}}^{\prime \prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{1}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{1}(\varphi, \vartheta)\right)^{T}+G_{f_{1}}^{\prime}\left(f_{1}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{1}(\varphi, \vartheta)\right\} q,\right. \\
& \left.G_{f_{2}}\left(f_{2}(\varphi, \ell)\right)-G_{f_{2}}\left(f_{2}(\varphi, \vartheta)\right)+\frac{1}{2} q^{T}\left\{G_{f_{2}}^{\prime \prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta} f_{2}(\varphi, \vartheta)\left(\nabla_{\vartheta} f_{2}(\varphi, \vartheta)\right)^{T}+G_{f_{2}}^{\prime}\left(f_{2}(\varphi, \vartheta)\right) \nabla_{\vartheta \vartheta} f_{2}(\varphi, \vartheta)\right\} q\right] .
\end{align*}
$$

Substituting the value of $f_{1}, f_{2}, G_{f_{1}}, G_{f_{2}}$ and $\eta_{2}$ at the point $\delta=0$ and simplify, we get

$$
\Pi_{3}=(4 v q \varphi,-1) \in-K \Rightarrow \Pi_{4}=\left(\ell^{4}+2 \varphi \ell^{2},-\ell\right) \in-K
$$

(A3). $\eta_{1}(\varphi, \delta)+\delta \in C_{1}, \forall \varphi \in C_{1}$.
(A4). $\eta_{2}(\ell, \vartheta)+\vartheta \in C_{2}, \forall \ell \in C_{2}$.
Validation: To prove our result its enough to prove that any point $\left(\varphi, 0, \lambda_{1}, \lambda_{2}, p\right)$ such that $\varphi \geqq 0, \lambda_{1}+\lambda_{2}=1$ are feasible to (GWPP). Also, the points $\left(0, \ell, \lambda_{1}, \lambda_{2}, q\right)$ such that $\ell \geqq 0, \lambda_{1}+\lambda_{2}=1$ are feasible to (GWDP). Now, at these feasible points,

$$
L=\left(L_{1}, L_{2}\right)=\left(\varphi^{2}-2 \varphi \lambda_{1} p^{2}, 1-2 \varphi \lambda_{1} p^{2}\right)
$$

and

$$
M=\left(M_{1}, M_{2}\right)=\left(\ell^{4}-\lambda_{1} q^{2}, 1-\ell-\lambda_{1} q^{2}\right)
$$

Now at above feasible condition

$$
\begin{equation*}
L-M=\left(\varphi^{2}-2 \varphi \lambda_{1} p^{2}-\ell^{4}+\lambda_{1} q^{2}, \ell-2 \varphi \lambda_{1} p^{2}+\lambda_{1} q^{2}\right) \notin K \backslash\{0\} \tag{28}
\end{equation*}
$$

In particular, the points $\left(\varphi, \vartheta, \lambda_{1}, \lambda_{2}, p\right)=\left(2,0, \frac{1}{2}, \frac{1}{2}, 1\right)$ and $\left(\delta, \ell, \lambda_{1}, \lambda_{2}, q\right)=$ $\left(0,2, \frac{1}{2}, \frac{1}{2}, 2\right)$ are
feasible for (GWPP) and (GWDP) respectively,
Now, calculate

$$
\begin{equation*}
L(\varphi, \vartheta, \lambda, p)-M(\delta, \ell, \lambda, q)=(-12,2) \notin K \backslash\{0\} . \tag{29}
\end{equation*}
$$

Hence, this validate the Weak duality Theorem 2.
Theorem 3 (Strong duality). Let $\left(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{p}_{1}=\bar{p}_{2}=\bar{p}_{3}=\ldots=\bar{p}_{k}\right.$ ) is an efficient solution of (GWPP); fix $\lambda=\bar{\lambda}$ in (GWDP) such that
(i) for all $i \in \tilde{N},\left[G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]$ is nonsingular,
(ii) the vector $\sum_{i=1}^{k} \bar{\lambda}_{i} \nabla_{\vartheta}\left[\bar{p}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{p}_{i}\right]$ $\notin \operatorname{span}\left\{G_{f_{1}}^{\prime}\left(f_{1}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{1}(\bar{\varphi}, \bar{\vartheta}), G_{f_{2}}^{\prime}\left(f_{2}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{2}(\bar{\varphi}, \bar{\vartheta}), \ldots, G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\}$,
(iii) the set of vectors $\left\{G_{f_{1}}^{\prime}\left(f_{1}(\bar{\phi}, \bar{\vartheta})\right) \nabla_{\theta} f_{1}(\bar{\varphi}, \bar{\vartheta}), G_{f_{2}}^{\prime}\left(f_{2}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\theta} f_{2}(\bar{\varphi}, \bar{\vartheta}), \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\theta} f_{k}(\bar{\varphi}, \bar{\vartheta})\right\}$ are linearly independent,
(iv) $\sum_{i=1}^{k} \bar{\lambda}_{i} \nabla_{\vartheta}\left[\bar{p}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{p}_{i}\right]=0 \Rightarrow \bar{p}_{i}=0, \forall i$, and
(v) $K$ is closed convex pointed cone with $R_{+}^{k} \subseteq K$.

Then, $\left(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{q}_{1}=\bar{q}_{2}=\bar{q}_{3}=\ldots=\bar{q}_{k}=0\right) \in W^{*}$ and $L(\bar{\varphi}, \bar{\vartheta}, \bar{p})=M(\bar{\varphi}, \bar{\vartheta}, \bar{q})$. Also, if the hypotheses of Theorem 1 or Theorem 2 are satisfied for all feasible solutions for (GWPP) and (GWDP), then $(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{p})$ and $(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{q})$ is an efficient solution for (GWPP) and (GWDP), respectively.

Proof. Since $\left(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{p}_{1}, \bar{p}_{2}, \bar{p}_{3}, \ldots ., \bar{p}_{k}\right)$, is an efficient solution of (GWPP), there exist $\alpha \in K^{*}, \beta \in C_{2}$ and $\bar{\eta} \in R$ such that the following Fritz -John optimality condition stated by [28] are satisfied at $\left(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{p}_{1}, \bar{p}_{2}, \bar{p}_{3}, \ldots, \bar{p}_{k}\right)$ :

$$
\begin{align*}
& (\varphi-\bar{\varphi})^{T}\left[\sum_{i=1}^{k} \alpha_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]+\sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta}) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]\left[\beta-\left(\bar{\alpha}^{T} e_{k}\right) \bar{\vartheta}\right]\right. \\
& \left.+\sum_{i=1}^{k} \bar{\lambda}_{i} \nabla_{\varphi}\left[\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{i}\right]\left(\beta-\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{\vartheta}+\frac{1}{2} \bar{p}_{i}\right)\right)\right] \geqq 0, \forall \varphi \in C_{1},  \tag{30}\\
& (\vartheta-\bar{\vartheta})^{T}\left\{\sum_{i=1}^{k} \alpha_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]+\sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]\right. \\
& \\
& \left(\bar{\beta}-\left(\bar{\alpha}^{T} e_{k}\right) \bar{\vartheta}\right)+\sum_{i=1}^{k} \bar{\lambda}_{i} \nabla_{\vartheta}\left[\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{i}\right] \\
& \quad\left[\bar{\beta}-\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{\vartheta}+\frac{1}{2} \bar{p}_{i}\right)\right]-\sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right.  \tag{31}\\
& +
\end{align*}
$$

$$
\begin{aligned}
& G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\bar{\beta}-\left(\bar{\alpha}^{T} e_{k}\right) \bar{\vartheta}\right)+\bar{\eta} e_{k}+\left\{\{ \beta - ( \overline { \alpha } ^ { T } e _ { k } ) ( \overline { \vartheta } + \frac { 1 } { 2 } \overline { p } _ { 1 } ) \} ^ { T } \left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}\right.\right. \\
& \left.+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{1},\left\{\{ \beta - ( \overline { \alpha } ^ { T } e _ { k } ) ( \overline { \vartheta } + \frac { 1 } { 2 } \overline { p } _ { 2 } ) \} ^ { T } \left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}\right.\right. \\
& \\
& \left\{\left\{\beta-\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{\vartheta}+\frac{1}{2} \bar{p}_{3}\right)\right\}^{T}\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{3}, \ldots,\right.
\end{aligned}
$$

$$
\begin{equation*}
\left\{\left\{\beta-\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{\vartheta}+\frac{1}{2} \bar{p}_{3}\right)\right\}^{T}\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{k}\right\}=0 \tag{32}
\end{equation*}
$$

$$
\begin{gather*}
{\left[G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]\left(\left(\bar{\beta}-\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{p}_{i}+\bar{\vartheta}\right)\right) \bar{\lambda}_{i}\right)=0, i \in \tilde{N}}  \tag{33}\\
\bar{\beta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{p}_{i}\right]=0  \tag{34}\\
\bar{\eta}^{T}\left[\bar{\lambda}^{T} e_{k}-1\right]=0  \tag{35}\\
(\bar{\alpha}, \bar{\beta}, \bar{\eta}) \geqq 0,(\bar{\alpha}, \bar{\beta}, \bar{\eta}) \neq 0 \tag{36}
\end{gather*}
$$

Inequalities (31) and (32) can be rewritten in the following expressions:

$$
\begin{gather*}
\sum_{i=1}^{k} \alpha_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]+\sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right] \\
\left(\bar{\beta}-\left(\bar{\alpha}^{T} e_{k}\right) \bar{\vartheta}\right)+\sum_{i=1}^{k} \bar{\lambda}_{i} \nabla_{\vartheta}\left[\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{i}\right] \\
{\left[\bar{\beta}-\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{\vartheta}+\frac{1}{2} \bar{p}_{i}\right)\right]-\sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right.} \\
\left.+\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{i}\right]\left(\bar{\alpha}^{T} e_{k}\right)=0 .  \tag{37}\\
G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\bar{\beta}-\left(\bar{\alpha}^{T} e_{k}\right) \bar{\vartheta}\right)+\left\{\left\{\beta-\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{\vartheta}+\frac{1}{2} \bar{p}_{i}\right)\right\}^{T}\right. \\
\left.\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{i}\right\}+\bar{\eta}=0, i \in \tilde{N} . \tag{38}
\end{gather*}
$$

Now, from hypothesis (iv), it is given that $R_{+}^{k} \subseteq K \Rightarrow \operatorname{int} K^{*} \subseteq \operatorname{int} R_{+}^{k}$.
Obviously, $\bar{\lambda}>0$ because $\bar{\lambda} \in \operatorname{int} K^{*}$.
By hypothesis (i), (33) gives

$$
\begin{equation*}
\beta=\left(\bar{\alpha}^{T} e_{k}\right)\left(\bar{p}_{i}+\bar{\vartheta}\right), i \in \tilde{N} . \tag{39}
\end{equation*}
$$

Suppose $\bar{\alpha}=0$, then (39) yields $\bar{\beta}=0$. Further, from (38) gives $\bar{\eta}=0$. Now, we reach at contradiction (36). Hence, $\bar{\alpha} \neq 0$. Further, $\bar{\alpha} \in K^{*} \subseteq R_{+}^{k}$ implies

$$
\begin{equation*}
\bar{\alpha}^{T} e_{k}>0 \tag{40}
\end{equation*}
$$

Now, we have to claim that $\bar{p}_{i}=0, \quad i \in \tilde{N}$. Using (39) and (40) in (38), we get

$$
\begin{array}{r}
\sum_{i=1}^{k} \bar{\lambda}_{i}\left[\nabla_{\vartheta}\left\{\frac{1}{2} \bar{p}_{i}\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{i}\right\}\right] \\
=-\frac{1}{\mu} \sum_{i=1}^{k}\left(\alpha_{i}-\mu \bar{\lambda}_{i}\right)\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right] \tag{41}
\end{array}
$$

By hypothesis (ii), we get

$$
\begin{equation*}
\sum_{i=1}^{k} \bar{\lambda}_{i}\left[\nabla_{\vartheta}\left\{\bar{p}_{i}\left(G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \bar{p}_{i}\right\}\right]=0 \tag{42}
\end{equation*}
$$

Again, from hypothesis (iv), we have

$$
\begin{equation*}
\bar{p}_{i}=0, \forall i \in \tilde{N} . \tag{43}
\end{equation*}
$$

From (39) implies

$$
\begin{equation*}
\bar{\beta}=\left(\bar{\alpha}^{T} e_{k}\right) \bar{\vartheta} . \tag{44}
\end{equation*}
$$

Using (42) and (43) in (37), we obtain

$$
\begin{equation*}
\sum_{i=1}^{k}\left(\alpha_{i}-\left(\bar{\alpha}^{T} e_{k}\right) \bar{\lambda}_{i}\right)\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]=0 \tag{45}
\end{equation*}
$$

From hypothesis (iii), it yields

$$
\begin{equation*}
\alpha_{i}=\left(\bar{\alpha}^{T} e_{k}\right) \bar{\lambda}_{i}, \quad i \in \tilde{N} . \tag{46}
\end{equation*}
$$

Using (43) and (44) in (30), we get

$$
\left.(\varphi-\bar{\varphi})^{T} \sum_{i=1}^{k} \alpha_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]\right] \geqq 0
$$

Using (40), (43), (44) and (46) in (30), we find that

$$
\begin{equation*}
\left.(\varphi-\bar{\varphi})^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]\right] \geqq 0, \quad \forall \varphi \in C_{1} . \tag{47}
\end{equation*}
$$

Let $\varphi \in C_{1}$. Then, $\varphi+\bar{\varphi} \in C_{1}$ and inequality (47) gives that

$$
\begin{equation*}
\left.\bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]\right] \geqq 0, \quad \forall \varphi \in C_{1} . \tag{48}
\end{equation*}
$$

Therefore,

$$
\begin{equation*}
\left.\sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]\right] \in C_{1}^{*} . \tag{49}
\end{equation*}
$$

Also, from (44), we obtain

$$
\begin{equation*}
\bar{\vartheta}=\frac{\bar{\beta}}{\bar{\alpha}^{T} e_{k}} \in C_{2} . \tag{50}
\end{equation*}
$$

Therefore, $\left(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{q}_{1}=\bar{q}_{2}=\bar{q}_{3}=\ldots=\bar{q}_{k}=0\right)$ satisfies the constraint of (GWDP) and is therefore a feasible solution for the dual problem (GWDP).

Now, letting $\varphi=0$ and $\varphi=2 \bar{\varphi}$ in (47), we obtain

$$
\begin{equation*}
\bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]=0 \tag{51}
\end{equation*}
$$

Further, from (34), (40), (43) and (44), we get

$$
\begin{equation*}
\bar{\vartheta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right]=0 \tag{52}
\end{equation*}
$$

Therefore, using (43), (51) and (52), we obtain

$$
\begin{aligned}
& \left(\begin{array}{l}
G_{f_{1}}\left(f_{1}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\vartheta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{p}_{i}\right] \\
\left.-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{p}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{p}_{i}\right], \ldots, G_{f_{k}}\left(f_{k}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\vartheta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right)\right. \\
\left.\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{p}_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{p}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}\right. \\
\\
\left.\left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta \vartheta \vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{p}_{i}\right]\right) \\
=\left(G_{f_{1}}\left(f_{1}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]\right. \\
\left.-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right], \ldots, G_{f_{k}}\left(f_{k}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right)\right. \\
\left.\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}\right. \\
\left.\left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]\right) .
\end{array}\right.
\end{aligned}
$$

This shows that the objective values are equal.
Finally, we have to claim that $\left(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{q}_{1}=\bar{q}_{2}=\bar{q}_{3}=\ldots=\bar{q}_{k}=0\right)$ is an efficient solution of (GWDP).

If possible, then suppose that $\left(\bar{\varphi}, \bar{\vartheta}, \bar{\lambda}, \bar{q}_{1}=\bar{q}_{2}=\bar{q}_{3}=\ldots=\bar{q}_{k}=0\right)$ is not an efficient solution of $(G W D P)$, then there exist $\left(\bar{\delta}, \bar{\ell}, \bar{\lambda}, \bar{q}_{1}=\bar{q}_{2}=\bar{q}_{3}=\ldots=\bar{q}_{k}=0\right)$ is efficient solution of (GWDP) such that

$$
\begin{aligned}
& \left(G_{f_{1}}\left(f_{1}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]\right. \\
& \left.-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right], \ldots, G_{f_{k}}\left(f_{k}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right)\right. \\
& \left.\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}\right. \\
& \\
& \left.\left.\quad+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\}_{\bar{q}}\right]-G_{f_{1}}\left(f_{1}(\bar{\delta}, \bar{\ell})\right)-\bar{\delta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}\right.\right. \\
& \left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right] \\
& \left.-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right], \ldots, G_{f_{k}}\left(f_{k}(\bar{\delta}, \bar{\ell})\right)-\bar{\delta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right)\right. \\
& \left.\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}\right. \\
& \\
& \left.\left.\left.+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right]\right) \in-K \backslash\{0\} .
\end{aligned}
$$

As

$$
\begin{aligned}
& \bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i} G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})=\bar{\vartheta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i} G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta}) \text { and } \bar{p}_{i}=0, \quad i \in \tilde{N}, \\
& \left(G_{f_{1}}\left(f_{1}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\vartheta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\vartheta} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]\right. \\
& \left.-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right], \ldots, G_{f_{k}}\left(f_{k}(\bar{\varphi}, \bar{\vartheta})\right)-\bar{\varphi}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right)\right. \\
& \left.\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\left(\nabla_{\varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right)^{T}\right. \\
& \left.\left.\quad+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\varphi}, \bar{\vartheta})\right\} \bar{q}_{i}\right]-G_{f_{1}}\left(f_{1}(\bar{\delta}, \bar{\ell})\right)-\bar{\delta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\varphi}, \bar{\vartheta})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})+\right. \\
& \begin{array}{r}
\left.\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right]
\end{array} \\
& \begin{array}{r}
\left.\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right], \ldots, G_{f_{k}}\left(f_{k}(\bar{\delta}, \bar{\ell})\right)-\bar{\delta}^{T} \sum_{i=1}^{k} \bar{\lambda}_{i}\left[G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right)\right.
\end{array} \\
& \left.\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})+\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right]-\frac{1}{2} \sum_{i=1}^{k} \lambda_{i} \bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}\right. \\
& \left.\left.\left.\quad+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right]\right) \in-K \backslash\{0\},
\end{aligned}
$$

which contradicts the Weak duality Theorem 1 or Theorem 2. Hence, completes the proof.
Theorem 4 (Converse duality). Let $(\bar{\delta}, \bar{\ell}, \bar{\lambda}, \bar{q})$ is an efficient solution of (GWDP); fix $\lambda=\bar{\lambda}$ in (GWPP) such that
(i) for all $i \in\{1,2, \ldots, k\},\left[G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right]$ is non singular,
(ii) $\sum_{i=1}^{k} \bar{\lambda}_{i} \nabla_{\varphi}\left[\bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right]$
$\notin \operatorname{span}\left\{G_{f_{1}}^{\prime}\left(f_{1}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{1}(\bar{\delta}, \bar{\ell}), G_{f_{2}}^{\prime}\left(f_{2}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{2}(\bar{\delta}, \bar{\ell}) \ldots, G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\}$.
(iii) the set of vectors $\left\{G_{f_{1}}^{\prime}\left(f_{1}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{1}(\bar{\delta}, \bar{\ell}), G_{f_{2}}^{\prime}\left(f_{2}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{2}(\bar{\delta}, \bar{\ell}), \ldots, G_{f_{k}}^{\prime}\left(f_{k}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{k}(\bar{\delta}, \bar{\ell})\right\}$ are linearly independent,
(iv) $\sum_{i=1}^{k} \bar{\lambda}_{i} \nabla_{\vartheta}\left[\bar{q}_{i}\left\{G_{f_{i}}^{\prime \prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\left(\nabla_{\varphi} f_{i}(\bar{\delta}, \bar{\ell})\right)^{T}+G_{f_{i}}^{\prime}\left(f_{i}(\bar{\delta}, \bar{\ell})\right) \nabla_{\varphi \varphi} f_{i}(\bar{\delta}, \bar{\ell})\right\} \bar{q}_{i}\right]=0 \Rightarrow \bar{q}_{i}=$ $0, \forall i$,
(v) $K$ is closed convex pointed cone with $R_{+}^{k} \subseteq K$.

Then, $(\bar{\delta}, \bar{\ell}, \bar{\lambda}, \bar{p}=0)$ is a feasible solution for $(G W P P)$ and the objective values of $(G W D P)$ and (GWPP) are equal. Furthermore, if the hypotheses of Theorem 1 or Theorem 2 are satisfied for all feasible solutions of $(G W D P)$ and $(G W P P)$, then $(\bar{\delta}, \bar{\ell}, \bar{\lambda}, \bar{p}=0)$ is an optimal solution of (GWPP). Also, if the hypotheses of Theorem 1 or Theorem 2 are satisfied for all feasible solutions for (GWDP) and (GWPP), then $(\bar{\delta}, \bar{\ell}, \bar{\lambda}, \bar{q})$ and $(\bar{\delta}, \bar{\ell}, \bar{\lambda}, \bar{p})$ is an efficient solution for (GWDP) and (GWPP), respectively.

Proof. It follows on the lines of Theorem 3.

## 4. Conclusions

In this paper, we have presented a novel generalized group of definitions and illustrated various non-trivial numerical examples for existing such type of functions. Numerical examples have also been illustrated to justify the weak duality theorem. Furthermore, we have studied a new class of $K-G_{f}$-Wolfe type primal-dual model with cone objective as well as constraint and proved duality theorem under $K-G_{f}$-bonvexity and $K-G_{f}$-pseudobonvexity. This work can further be extended to higher order symmetric
fractional programming problem and variational control problem over cones. This will be feature task for the researchers.

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## Article

# Predicting Sit-to-Stand Body Adaptation Using a Simple Model 

Sarra Gismelseed * ${ }^{(D)}$, Amur Al-Yahmedi © , Riadh Zaier, Hassen Ouakad (D) and Issam Bahadur (D)<br>Department of Mechanical \& Industrial Engineering, College of Engineering, Sultan Qaboos University, Muscat 123, Oman; amery@squ.edu.om (A.A.-Y.); zaier@squ.edu.om (R.Z.); houakad@squ.edu.om (H.O.); bahdoor@squ.edu.om (I.B.)<br>* Correspondence: sarraabbasher@gmail.com

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[^1]
#### Abstract

Mathematical models that simulate human motion are used widely due to their potential in predicting basic characteristics of human motion. These models have been involved in investigating various aspects of gait and human-related tasks, especially walking and running. This study uses a simple model to study the impact of different factors on sit-to-stand motion through the formulation of an optimization problem that aims at minimizing joint torques. The simulated results validated experimental results reported in the literature and showed the ability of the model to predict the changes in kinetic and kinematic parameters as adaptation to any change in the speed of motion, reduction in the joint strength, and change in the seat height. The model discovered that changing one of these determinants would affect joint angular displacement, joint torques, joint angular velocities, center of mass position, and ground reaction force.


Keywords: sit-to-stand; discrete mechanics; optimization; sit-to-stand determinates

MSC: 37N99

## 1. Introduction

Mathematical models that simulate human motion and capture the basic pattern of human motor characteristics in the absence of empirical data (or what is known as predictive simulation [1]), were used to study many human kinetic characteristics, such as targeting a specific kinetic speed while reducing the metabolic cost associated with this movement [1,2]. By applying dynamic optimization processes while considering physical and physiological constraints [3,4], researchers were able to derive whole body dynamics, aiding in the production of movements that closely resemble human motion by controlling relevant parameters. Most of the predictive simulation studies were done to study human gait characteristics since it is the main feature of human beings [3,5,6]; however, Sit-to-Stand (STS) motion is also one of the main daily activities, which is also considered as a distinctive feature of human beings.

STS motion has been attracting scientists and physical therapists to experimentally study and analyze the aspects of how humans perform this task that is mainly defined by the process of going to a standing position from a sitting position [7,8]. This motion has a direct effect on humans' quality of life, and the significance of this motion will not be appreciated until it becomes physically or cognitively challenging, which can be temporary due to injury or permanent due to aging or illness. All the previous reasons may cause difficulties in completing this motion successfully, especially in an elderly population. Studies showed that at the age of 55 years and above, $45 \%$ of women and $30 \%$ of men suffer from moderate to serious inabilities in rising from sitting [9]. Inability to properly perform the STS motion may be associated with decreased mobility and balance, with an increased risk of falling, as well. Several factors influencing the performance of STS motion are considered as determinants, as they can either facilitate successful completion of the task or make it more challenging [10,11]. The determinants of STS motion have
been categorized into different groups according to the source of the determinates, which are strategy-related determinants such as speed, foot position, trunk rotation, etc. [8,11]; subject-related determinants such as age, muscle force, and disease [7,12,13]; and chairrelated determinants $[14,15]$ such as seat height, armrests, and backrest $[16,17]$. The impact of these determinants on body behavior and movement performance differs according to the type of determinants, and it has been studied in order to investigate how the body adapts to the changes of these determinants.

The predictive simulation has been involved in studying the STS motion $[6,18]$ to have a better understanding of how we usually move and when this movement occurs [18]. A good understanding of the normal way of movement during STS will help as well in applying this knowledge to some cases of predictive simulation that involve studying losing mobility or some difficulties in performing this task. The complexity of biomechanical models used to study STS motion varies from simple to complicated based on the number of segments and muscles included in these models, in addition to the number of planes used in the modeling [9,19]. Most of the models were modeled as 2D models in the sagittal plane [20] since most of the joints' movement occur in this plane; however, few were modeled as 3D models. These models were involved in studying various aspects of STS motion and the effect of different parameters on this motion, including the effect of chair height, the effect of initial foot position, the effect of muscle strength, etc. Garner proposed a biomechanical model modeled in the sagittal plane that has three rigid segments (i.e., thigh, leg, and HAT (Head, Arms, and Torso) and eight muscles [19]. Garner used dynamic optimization to explain the execution of this motion by minimizing a cost function split into two parts; the first part is aimed at minimizing muscle stress until the time the biped leaves the chair, and the second part is aimed at minimizing peak forces developed by the muscles from the time the biped leaves the chair until the time required to accomplish this task [19]. This cost function produced simulation results that were very similar to the experimental results. The model of Garner is further enhanced by Daigle by adding a foot segment and increasing the number of muscles to 18 [20], and the model was used to understand the difficulties related to muscle strength and its effect on accomplishing this task by using an optimization formulation that minimizes the motion time with different muscle strength varied from $50 \%$ to $200 \%$ [20]. The model was able to define which muscles have a great contribution and power activation in the sit-to-stand task [20]. On other hand, Domire used a sagittal plane model that consists of three links connected with three joints and actuated by eight muscles to study the effect of seat height on STS motion [21]. Domire predicted the effect of seat height using the same objective function described by Garner and simulated the movement during STS from different seat heights, adjusted by rotating the thigh segment at four different angles: $80^{\circ}, 90^{\circ}, 98^{\circ}$, and $100^{\circ}$ [21]. The study concluded that as the seat height decreased, the movement of STS became more difficult, and it was impossible at the lowest seat height corresponding to the $100^{\circ}$ thigh angle.

In this study, a simple 2D model is used to predict the effect of changing three STS determinants (speed, joint strength, and seat height) on lower limb and upper body kinetic and kinematic parameters by applying dynamic optimization to minimize the joint torques. First, the derivation of the proposed model and the formulation of the optimization problem are shown in the following section. Then, the simulation results are shown in the third section of this study with the validation of some simulation results with experimental results from the literature. In the fourth section, results from this study are discussed and compared to reported results and conclusions from experimental studies. Finally, we summarize our work with the main conclusions in the last section.

## 2. Materials and Methods

Accurately predicting STS motion requires a biomechanics model that accurately represents the human body, as well as an understanding of how to identify the specific movements involved in STS motion. The model proposed in this study is modeled in the sagittal plane since the movement of humans' joints are mostly obvious when it is seen
from the sagittal plane compared to the other planes [22]. This will allow for simplifying the human body to a three-links model, as shown in Figure 1, through assuming that the two legs are symmetrical and can be modeled as one leg with two segments, tibia and femur, in addition to the torso segment. A male subject of height 1.70 m and mass 70 kg is considered for approximating the model physical parameters. The physical parameters shown in Table 1 were calculated using anthropometric percentages compiled by different investigators (the detailed calculations are given in Appendix A) [23-25].


Figure 1. Three-links model, $u_{i}$ is the torque applied at each joint, and each link is assigned with a number: $1=$ stance tibia, $2=$ stance femur, $3=$ trunk.

Table 1. Physical parameters.

| Link | Body <br> Segment | Length of <br> Segment <br> $(\mathbf{m})$ | Distance from <br> Center of Mass <br> to Next Joint (m) | Mass of <br> Segment <br> $(\mathbf{k g})$ | Moment of Inertia <br> of Segment <br> $\left(\mathbf{k g} \cdot \mathbf{m}^{\mathbf{2}} \mathbf{)}\right.$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $\boldsymbol{l}_{\boldsymbol{i}}$ | 0.390 | 0.182 | $\boldsymbol{m}_{\boldsymbol{i}}$ |

The problem is formulated as an optimization problem in the goal is to minimize a given cost function while satisfying constraints imposed by the task. Due to the high degree of nonlinearity of the problem, the use of the direct (single) shooting method to solve the resulting optimization problem may likely cause the algorithm to fail to find a solution. One can either use collocation methods or use direct discretization of the Lagrange-d'Alembert Principle for the system. We chose to use direct discretization of the Lagrange-d'Alembert Principle (discrete mechanics) [26,27]. Discrete mechanics requires that the Lagrange equation of the system be derived first. The Lagrange equation of the model is given by:

$$
\begin{align*}
\mathcal{L}=\frac{1}{2}\left(m_{1}\left(l_{1}-d_{1}\right)^{2}\right. & \left.+m_{2} l_{1}^{2}+m_{3} l_{1}^{2}\right) \dot{\theta}_{1}^{2}+\frac{1}{2}\left(m_{2}\left(l_{2}-d_{2}\right)^{2}+m_{3} l_{2}^{2}\right) \dot{\theta}_{2}^{2}+\frac{1}{2} m_{3} d_{3}^{2} \dot{\theta}_{3}^{2} \\
& -\left(m_{2} l_{1}\left(l_{2}-d_{2}\right)+m_{3} l_{1} l_{2}\right) \dot{\theta}_{1} \dot{\theta}_{2} \cos \left(\theta_{1}+\theta_{2}\right)+m_{3} l_{1} d_{3} \dot{\theta}_{1} \dot{\theta}_{3} \cos \left(\theta_{1}-\theta_{3}\right)  \tag{1}\\
& -m_{3} l_{2} d_{3} \dot{\theta}_{2} \dot{\theta}_{3} \cos \left(\theta_{2}+\theta_{3}\right)-g\left(m_{1}\left(l_{1}-d_{1}\right)+m_{2} l_{1}+m_{3} l_{1}\right) \cos \theta_{1} \\
& -g\left(m_{2}\left(l_{2}-d_{2}\right)+m_{3} l_{2}\right) \cos \theta_{2}-m_{3} g d_{3} \cos \theta_{3}
\end{align*}
$$

the discrete Lagrange equation is derived from the continuous Lagrange equation, using the mid-point rule as follows (for the detailed derivations, check Appendix A):

$$
\begin{equation*}
\mathcal{L}(q, \dot{q}) \rightarrow h \mathcal{L}_{d}\left(\frac{q_{k+1}+q_{k}}{2}, \frac{q_{k+1}-q_{k}}{h}\right) \tag{2}
\end{equation*}
$$

The system dynamics is then:

$$
\begin{align*}
D_{1} \mathcal{L}_{d}\left(\theta_{1 k}, \theta_{1 k+1},\right. & \left.\theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)+f_{d}^{-}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& +D_{2} \mathcal{L}_{d}\left(\theta_{1 k-1}, \theta_{1 k}, \theta_{2 k-1}, \theta_{2 k}, \theta_{3 k-1}, \theta_{3 k}\right)  \tag{3}\\
& +f_{d}{ }^{+}\left(\theta_{1 k-1}, \theta_{1 k}, \theta_{2 k-1}, \theta_{2 k}, \theta_{3 k-1}, \theta_{3 k}\right)=0
\end{align*}
$$

where $D_{1} \mathcal{L}_{d}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)$ is the first derivative of the discrete Lagrange with respect to current coordinates (i.e., $\theta_{1 k}, \theta_{2 k}$, and $\theta_{3 k}$ ), and is the first derivative of the discrete Lagrange with respect to future coordinates (i.e., $\theta_{1 k+1}, \theta_{2 k+1}$, and $\theta_{3 k+1}$ ).

$$
\begin{gather*}
D_{1} \mathcal{L}_{d}=\left[\begin{array}{l}
\frac{\partial \mathcal{L}_{d}}{\partial \theta_{1 k}}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
\frac{\partial \mathcal{L}_{d}}{\partial \theta_{2 k}}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
\frac{\partial \mathcal{L}_{d}}{\partial \theta_{3 k}}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)
\end{array}\right]  \tag{4}\\
D_{2} \mathcal{L}_{d}=\left[\begin{array}{l}
\frac{\partial \mathcal{L}_{d}}{\partial \theta_{1 k+1}}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
\frac{\partial \mathcal{L}_{d}}{\partial \theta_{2 k+1}}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
\frac{\partial \mathcal{L}_{d}}{\partial \theta_{3 k+1}}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)
\end{array}\right] \tag{5}
\end{gather*}
$$

Both discrete torques, i.e., left and right forces, are defined as follows.

$$
u_{d}^{-}=u_{d}^{+}=\left[\begin{array}{c}
u_{1 k}-u_{2 k}  \tag{6}\\
u_{2 k}-u_{3 k} \\
u_{3 k}
\end{array}\right]
$$

To raise the body from the sitting position to the standing position, it is necessary to provide an appropriate amount of torque at the joints. Therefore, an optimization problem is formulated to mimic the Sit-to-Stand function with the aim of minimizing the effort required to complete this motion, which is represented as the total torque-squared applied by each joint [28]. The torques are equally minimized ( $\alpha_{1}=\alpha_{2}=\alpha_{3}=1$ ) to study the effect of speed and seat height on the motion; however, $\alpha_{2}$ (coefficient of the knee torque) is increased to study the effect of knee strength while keeping $\alpha_{1}=\alpha_{3}=1$, assuming that the model is mimicking sitting on a desk chair in which the femur is making a $90^{\circ}$ angle with the vertical axis, whereas the tibias and torso are allowed to have unfixed initial states but within reasonable bounded range. Conversely, the standing position requires all joints to make around $0^{\circ}$ with the vertical as defined in our model. Then, the optimization problem can be simplified as follows:

$$
\begin{equation*}
\min J=\sum_{k=1}^{N} \sum_{i=1}^{3} \alpha_{i} u_{i}^{2}(N)=\sum_{k=1}^{N} \alpha_{1} u_{1}^{2}(N)+\alpha_{2} u_{2}^{2}(N)+\alpha_{3} u_{3}^{2}(N) \tag{7}
\end{equation*}
$$

with both decision variables $\theta_{i}$ and $u_{i}$ subjected to the following constraints:

- System dynamics

$$
\begin{aligned}
D_{1} \mathcal{L}_{d}\left(\theta_{1 k}, \theta_{1 k+1},\right. & \left.\theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)+f_{d}^{-}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& +D_{2} \mathcal{L}_{d}\left(\theta_{1 k-1}, \theta_{1 k}, \theta_{2 k-1}, \theta_{2 k}, \theta_{3 k-1}, \theta_{3 k}\right) \\
& +f_{d}{ }^{+}\left(\theta_{1 k-1}, \theta_{1 k}, \theta_{2 k-1}, \theta_{2 k}, \theta_{3 k-1}, \theta_{3 k}\right)=0
\end{aligned}
$$

- Boundary constraints

$$
\theta_{3}(1)=0^{\circ}, \theta_{2}(1)=90^{\circ}
$$

$$
\theta_{1}(N+1)=0^{\circ}, \theta_{2}(N+1)=0^{\circ}, \theta_{3}(N+1)=0^{\circ}
$$

- Path constraints

$$
\begin{aligned}
\theta_{i \min } & \leq \theta_{i} \leq \theta_{i \max } \\
u_{i \min } & \leq u_{i} \leq u_{i \max }
\end{aligned}
$$

The optimization problem was solved using the MATLAB ${ }^{\circledR}$ and SNOPT $^{\circledR}$ (Sparse Nonlinear OPTimizer) toolbox, and the simulation was run with different initial conditions for each case in this study.

## 3. Results

The results of the optimization provided a unique solution for each case that describes the optimal STS movement for the given conditions. The normal pattern of STS can be explained using the ground force and the velocity of Center of Mass (COM) by dividing the STS motion into phases, starting from the model at sitting position and ending when the velocity of center of mass approaches zero [29]. The first phase starts at the sitting position and initiates the movement by bending the torso forward, and it ends by raising both thighs, as seen in Figure 2. Bending the torso forward helps to get torso flexion, which shifts COM forward. This phase transfers momentum to the next phase, and it occurs quickly. Therefore, it requires the generation of an efficient flexion momentum to be able to transfer a lot of it to the lower extremities in the next phase. Before standing up, the tibia segment is placed backward behind the knee, which results in ankle dorsiflexion. The results showed that the profile of GRF directed forward under the feet and the velocity of COM have a bell shape with a peak just before seat off, as seen in Figures 3 and 4 respectively [29]. The GRF for the model and the experimental results [30] have the same pattern; however, for the case of the experimental results, the GRF starts from zero due to the presence of the chair that holds the weight of the person. For the model, the GRF as a percentage of body weight starts from one because we assume that the model is at the sitting position on foot (no chair). Moreover, we noticed that the model obtained the same profile of COM velocity reported in the literature [29].


Figure 2. The normal pattern of STS can be explained using the ground force reaction and the velocity of center of mass by dividing the movement into phases starting from the model sitting and ending when the velocity of center of mass approaches zero.

Then, the model is considered to be in the second phase once it is no longer in the sitting position. The hip will continue to flex, and the model uses the flexion momentum transferred from the previous phase, which was generated from torso flexion and distributed to all segments, but more to the lower segments. For human beings, the flexion momentum produced in the first phase is for the upper body, since the thighs are still on the seat. Whenever the thighs are taken off the seat, this flexion momentum can be transferred to the total body and into the legs, as well. Meanwhile, we will have continued flexion of the hips and continued ankle dorsiflexion in this phase. One of the most important
characteristics of this phase is that it has the largest amount of ground reaction force among all four phases of STS. This is because during this phase, all the body weight is on the feet. Looking at the pattern of Ground Reaction Force (GRF), we see that there is a basic need to produce ground force greater than body weight in order to accelerate the body in the upward direction [30].


Figure 3. The profile of GRF directed forward under the feet from simulation and experimental results [30].


Figure 4. The velocity of COM of the model and the experimental results [29].
The third phase starts at maximum ankle dorsiflexion and ends at hip extension. It includes a sequence of lower limb extensions, knee extension, hip extension, and ankle extension. When a person stands up to a standing position from the point of maximum ankle dorsiflexion, the foot does not change position, but the tibia will move backward on the ankle. Therefore, the angle between tibia and the foot will be around $60^{\circ}$ in this phase, whereas in the standing position it will be around $90^{\circ}$. In the previous phase, muscles are activated to flex the knee, whereas in this phase, the muscles are activated to extend the knee.

The GRF will start to decrease from its maximum value, indicating that the thigh is leaving the seat [30]. Again, the ground reaction force will start to increase once the body is in the upward standing position [31]. This phase is terminated once the hip reaches its maximum extension. The last phase is mainly about stabilizing the body and preventing
it from falling down, since the body is already in a standing position. In this phase, the velocity of COM will approach zero, indicating the termination of STS movement.

### 3.1. STS Speed

The effect of STS speed on joints' kinetic and kinematics was investigated using the three-links model by changing the time required to stand up from sitting position. In order to stand up from the sitting position, the model flexes the torso forward with the presence of the hip flexion. The results show that as the speed increases (time to stand up is reduced), the torso remains flexed for a longer time interval as a percentage of total STS duration, as seen in Figure 5. Additionally, at very short time of STS, the model stands up with a flexed torso. This means that the duration of the stabilization phase reduces as the speed increases. In addition to the increment in the flexion duration of the torso, the torso angular velocity tends to increase as the speed increases, as shown in Figure 6.


Figure 5. (a) The model performing the sit-to-stand task within different timeframes; (b) Torso angle during standing from sitting position within different timeframes.


Figure 6. Torso angular velocity at different STS speeds.

An initial flexion angular velocity (negative) is observed at the hip joint as the torso is leaning forwards, followed by a positive angular velocity indicating the extension of the hip joint, which is continuous until the standing position is attained, as is clear from Figure 7. On the other hand, the velocity of the knee joint shows a smooth increase and decrease in the negative direction (extension velocity), demonstrating the control of the movement towards the full extension position, as shown in Figure 7. The peak of the angular velocity increases as the speed of rising increases, indicating the presence of exaggerated knee joint loading and causing an increment in the knee joint torque, as portrayed in Figure 8.


Figure 7. Hip joint velocity and knee joint velocity at different STS speeds.


Figure 8. Ankle, knee, and hip torques at different STS speeds.
A sufficient torque needs to be generated at the joints to overcome the challenges that will need to be faced to accomplish this task successfully, including moving the body's center of mass forward and raising it vertically to the standing position.

The results have shown that the increment in the speed of STS increases the ankle dorsiflexion, knee extension, and hip flexion joint torques, and among the models' joints, the results of the optimization demonstrated that the knee is the most affected joint by increasing speed, as displayed in Figure 8.

### 3.2. Reduction in Joint Strength

To investigate the effect of reduced knee strength on STS motion, the model is used to optimize the STS motion with minimum joint torques and especially the knee joint torque. This can be accomplished by reformulating the cost function in Equation (7) to have a larger
coefficient for the knee joint torque. The optimizer was run with different coefficients for the knee joint torque ( $\alpha_{2}=1,3,5$, and 7). The results of the optimization showed that as the coefficient of knee joint torque increases, the model tends to flex the knee and the torso as well while standing up and start to extend them at the moment just before the standing position, which is as shown in Figure 9.


Figure 9. The model standing up with reduced knee joint strength corresponds to different coefficients of knee joint torque ( $\alpha_{2}=1, \alpha_{2}=3, \alpha_{2}=5$, and $\alpha_{2}=7$ ).

Increasing the coefficient of knee joint torque to indicate the difficulties in rising due to the knee joint seems to increase the angular velocity of hip flexion and angular velocity of knee extension as seen in Figure 10. Moreover, an increment in the hip joint torque was observed in Figure 11 when increasing the coefficient of knee joint torque, indicating that the model depends on the hip joint to accommodate for the losing of some knee joint strength.


Figure 10. The angular velocity of the model with reduced knee joint strength corresponds to different coefficients of knee joint torque ( $\alpha_{2}=1, \alpha_{2}=3, \alpha_{2}=5$, and $\alpha_{2}=7$ ).


Figure 11. The Hip and knee joint torques of the model with reduced knee joint strength correspond to different coefficients of knee joint torque ( $\alpha_{2}=1, \alpha_{2}=3, \alpha_{2}=5$, and $\alpha_{2}=7$ ).

Furthermore, the effect of reduced knee strength was studied at two speeds, fast and slow, corresponding to STS duration of 1.5 s and 0.8 s , respectively. The results of different STS speeds showed a similar pattern of joint motion, as shown in Figure 12. The model first flexed the hip while still at the sitting position with the knee initially flexed at $90^{\circ}$, and then performed a series of extension movements by extending the hip and knee (after a short delay) joints while plantar flexing the ankle during the remaining time of the standing. Immoderate hip flexion is resisted by the contraction of muscles, which at the same time motivates the knee flexion just prior standing up. Knee flexion is also controlled by the contraction of muscles to avoid excessive knee flexion. Figure 12 shows that the joint angle profiles for the case of larger knee torque coefficients were identical to the normal case at a slow speed; however, the range of joint motion with reduced knee strength was slightly decreased at a fast speed. On the other hand, for the same case (normal knee strength or reduced knee strength), increasing the speed resulted in increased joint range of motion.


Figure 12. Hip angle, Knee angle, and Ankle angle at different STS speeds with normal and reduced knee joint strength ( $\alpha_{2}=1$ and $\alpha_{2}=5$, respectively).

With normal knee strength, and as a consequence of a more upright torso during standing up, as stated in the previous section, the body's center of mass exhibited a more posterior location (behind the foot) in the cases of a fast speed that corresponds to a standing duration of 0.8 s and 1.0 s . However, with reduced knee strength, the increment in the speed increases the horizontal position of the center of mass, as seen in Figure 13. On the
other hand, with normal knee strength and reduced knee strength, the vertical position of the body's center of mass starts to increase earlier for the cases of slow motion, indicating that the torso is remaining flexed while standing up for the cases of fast motion. However, the impact of speed on the vertical position of center of mass is clearer for the case of reduced knee strength than normal knee strength.


Figure 13. Horizontal and vertical positions of the center of mass at different STS speeds with normal and reduced knee joint strength ( $\alpha_{2}=1$ and $\alpha_{2}=5$, respectively).

Figure 14 shows that the increment and decrement in the angular velocity of the hip and knee joints tends to be rough at a fast STS motion for $\alpha_{2}=5$ with a significant increment in the peak value of knee angular velocity compared to the hip angular velocity.


Figure 14. Hip joint velocity and knee joint velocity at different STS speeds with normal and reduced knee joint strength ( $\alpha_{2}=1$ and $\alpha_{2}=5$, respectively).

### 3.3. Seat Height

The results of the optimization of four seat heights ( $30 \mathrm{~cm}, 33 \mathrm{~cm}, 37 \mathrm{~cm}$, and 44 cm ) showed almost a similar time taken by the model to stand from the sitting position, although the time required to attain the full standing position was longer for the lowest seat height. Moreover, the torso angle showed noticeable differences in the forward movement by showing an increment in the maximum torso angle, as seen in Figure 15, as the seat height decreases with larger angular displacements of the joints due to the larger space needing to be covered.


Figure 15. (a) Model standing up from different seat heights; (b) Torso angle of the model standing up from four seat heights ( $30 \mathrm{~cm}, 33 \mathrm{~cm}, 37 \mathrm{~cm}$, and 44 cm ).

The results apparently showed that as the seat height decreases, the hip and knee extension velocity increases, as seen in Figure 16, resulting in more foot repositioning, which led to an increase in the ground reaction force, as shown in Figure 17.


Figure 16. Hip and knee joint velocities of the model standing up from four seat heights $(30 \mathrm{~cm}$, $33 \mathrm{~cm}, 37 \mathrm{~cm}$, and 44 cm ).


Figure 17. Ground Reaction Force of the model (as a percentage of body weight) standing up from four seat heights ( $30 \mathrm{~cm}, 33 \mathrm{~cm}, 37 \mathrm{~cm}$, and 44 cm ).

We have found that lowering the seat height increases the hip torque before and after the seat off, as well; however, reducing the height decreases the knee torque before seat off while increases it after seat off, as clear from Figure 18.


Figure 18. Hip joint and knee joint torque of the model standing up from four seat heights ( 30 cm , $33 \mathrm{~cm}, 37 \mathrm{~cm}$, and 44 cm ).

## 4. Discussion

We have found that increasing STS speed resulted in a shorter duration of STS phases, which agrees with reported experimental studies in the literature $[32,33]$ and a more upright position of the torso at lift-off. Furthermore, in order to speed up the motion of STS by reducing the duration of phases, the joint torque was significantly affected, which agrees with reported experimental studies in the literature $[7,13,34,35]$. It is also found that increasing the speed of STS motion increased the hip flexion torque, knee extension torque, and ankle dorsiflexion torque for the model with both normal knee strength and reduced knee strength, and this validates the results of other experimental studies on healthy people and people with weakened knees [7,13,34,35]. However, between the model joints, the knee joint was the most affected joint, and it was clear from the significant increment in the knee extension moment and knee extension velocity [8,35]. This impact may be due to the load of body weight on the knee at the moment of standing up, since the model was standing with a flexed torso.

Therefore, due to the great increment in the knee joint torque, more attention was paid to the knee joint by increasing the coefficient of knee joint torque in the cost function to represent individuals with weakened or impaired knee joints, such as patients with some knee osteoarthritis or elderly people, and investigated the effect of this on STS motion. Furthermore, it was observed that as the strength of the knee joint was reduced (increasing the coefficient of knee joint torque), the angular velocities of hip and knee joints slightly increased; however, the increment in the hip joint torque was very significant. This indicated that the hip joint was trying to compensate for the reduction in the available knee joint torque. To avoid the load on the knee joint, a person can stand up slowly, especially when there are abnormalities in the knee joint; however, fast STS motion is required not only as a daily life activity, but also to do some clinical settings and tests like fast STS tests, in which patients are asked to stand up as quickly as possible. Regarding the kinetic parameters of lower limb joints and how they are affected by speed, the results of fast STS motion for the model with normal and reduced knee strength was characterized by increased knee and hip peak velocities, similar results being reported in experimental studies in the literature $[7,8]$ and they tended to be unsmooth at fast speed as we increased the coefficient of knee joint torque. We have found that as the speed of motion was increased, the horizontal position of center of mass also increased, which can be related to the upright position of the
torso while standing up, and the results of this simulation are in agreement with reported experimental studies reported in the literature [35,36]. However, the pattern of lower limb joints, including the hip, was similar when performing STS motion at fast and slow speeds for both normal knee strength and reduced knee strength [36].

Conversely, we have investigated the effects of seat height on the kinetic and kinematic parameters of the model during STS motion, as it has been considered as the most influential factor related to the chair properties $[16,37,38]$. The results indicated that reducing the seat height increases the angular velocity of the hip and knee joints required to stand up, which was a consequence of the increment in the torso angular velocity required to move the torso further forward from a lower seat height before seat off $[16,37,38]$. The higher angular velocity resulted in a higher torque at the hip joint, which transferred to the lower limb joints after leaving the seat. It is also noticed that the model exhibited an increment in the torso movement which was demonstrated by Weiner et al. as a result of increasing the demand to move the center of mass closer to the knees and reduce the required effort [39]. The higher torque at the hip joint was used by the model to assist in the remaining phases of STS, which was also due to the high forward velocity of the torso at seat off [39]. The increment in the joints' angular displacement exhibited by the model as the seat height is decreased is also confirmed by other studies [11,40], and the reason behind that was demonstrated by Yamada and Demura as reported in [17] that standing up from a lower seat height is more difficult due to the increment in the distance that must be covered by the person in order to stand up. The increase in the distance covered will also lead to an increment in the muscle activity of the lower limbs, which is represented by the increment of the model's joint torques in this study [11]. It has been known that the increment in the torso's forward movement increases the chance of falling down due to the huge torso mass compared to lower limbs; therefore, more feet repositioning was observed by the model when lowering the seat height as a stabilizing strategy $[37,38]$, which resulted in an increment in the peak of the ground reaction force.

## 5. Conclusions

In this study, we have used a mathematical model to investigate its ability to predict STS body adaptations by formulating an optimization problem that minimizes joints' torques. The model was able to predict basic features of STS motion by following the constraints of the optimization problem. We have found that STS motion may be influenced by different factors such as motion speed, reduced joint strength, and seat height. We have first investigated the effect of motion speed by changing the total time required to stand up from the sitting position while equally optimizing the joint torques. Then, we reduced the knee joint strength, since it was the most affected joint, by increasing the coefficient of joint torque in the cost function for the purpose of optimizing it more than the other joints torques. Finally, we studied the effect of seat height on the STS motion by changing the femur link to get different seat heights, and again equally optimized the joints torque.

The results of this study for the three cases agreed with published experimental results, indicating that the model was able to predict the STS body adaptation. Increasing STS motion speed, reducing knee strength, and reducing seat height was found to increase joints torque and joints angular velocity, whereas reducing the knee strength was found to decrease the range of joints motion. On the other hand, we have found that reducing knee strength and reducing seat height led to increased torso flexion, while increasing STS motion speed led to standing up with an upright torso position. Moreover, the model expressed higher ground reaction force as the height of the seat was reduced and more posterior center of mass position as the knee strength was reduced.

Since STS motion is used for clinical investigation as well as daily activity as human beings, it is very important to pay attention to the factors that may impact this motion. Therefore, the model involved in this study can help in understanding these factors and determining the ability of the person to complete this motion and the challenges that may be faced.

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## Appendix A

The model under study consists of three segments: the torso, the femur, and the tibia. This model is simplified as we study human motion in the sagittal plane in which we can assume that the two legs move in an asymmetric manner, so we can include one leg only. The system dynamics are derived using direct discretization of the Lagrange-d'Alembert Principle (discrete mechanics). First, we define the positions of center of mass of each link in the $x$ and $y$ coordinates:

$$
\begin{gather*}
x_{c 1}=-\left(l_{1}-d_{1}\right) \sin \theta_{1}  \tag{A1}\\
y_{c 1}=\left(l_{1}-d_{1}\right) \cos \theta_{1}  \tag{A2}\\
x_{c 2}=-l_{1} \sin \theta_{1}-\left(l_{2}-d_{2}\right) \sin \theta_{2}  \tag{A3}\\
y_{c 2}=l_{1} \cos \theta_{1}+\left(l_{2}-d_{2}\right) \cos \theta_{2}  \tag{A4}\\
x_{c 3}=-l_{1} \sin \theta_{1}-l_{2} \sin \theta_{2}-\left(l_{3}-d_{3}\right) \sin \theta_{3}  \tag{A5}\\
y_{c 3}=l_{1} \cos \theta_{1}+l_{2} \cos \theta_{2}+\left(l_{3}-d_{3}\right) \cos \theta_{3} \tag{A6}
\end{gather*}
$$

The velocities of center of mass of each link in the $x$ and $y$ coordinates are:

$$
\begin{gather*}
\dot{x}_{c 1}=-\left(l_{1}-d_{1}\right) \dot{\theta}_{1} \cos \theta_{1}  \tag{A7}\\
\dot{y}_{c 1}=-\left(l_{1}-d_{1}\right) \dot{\theta}_{1} \sin \theta_{1}  \tag{A8}\\
\dot{x}_{c 2}=-l_{1} \dot{\theta}_{1} \cos \theta_{1}-\left(l_{2}-d_{2}\right) \dot{\theta}_{2} \cos \theta_{2}  \tag{A9}\\
\dot{y}_{c 2}=-l_{1} \dot{\theta}_{1} \sin \theta_{1}-\left(l_{2}-d_{2}\right) \dot{\theta}_{2} \sin \theta_{2}  \tag{A10}\\
\dot{x}_{c 3}=-l_{1} \dot{\theta}_{1} \cos \theta_{1}-l_{2} \dot{\theta}_{2} \cos \theta_{2}-d_{3} \dot{\theta}_{3} \cos \theta_{3}  \tag{A11}\\
\dot{y}_{c 3}=-l_{1} \dot{\theta}_{1} \sin \theta_{1}-l_{2} \dot{\theta}_{2} \sin \theta_{2}-d_{3} \dot{\theta}_{3} \sin \theta_{3} \tag{A12}
\end{gather*}
$$

The kinetic energy of each link:

$$
\begin{equation*}
K E_{1}=\frac{1}{2} m_{1} v_{c 1}^{2}=\frac{1}{2} m_{1}\left(l_{1}-d_{1}\right)^{2} \dot{\theta}_{1}^{2}+\frac{1}{2} I_{1} \dot{\theta}_{1}^{2} \tag{A13}
\end{equation*}
$$

$$
\begin{equation*}
K E_{2}=\frac{1}{2} m_{2} v_{c 2}^{2}=\frac{1}{2} m_{2} l_{1}^{2} \dot{\theta}_{1}^{2}+\frac{1}{2} m_{2}\left(l_{2}-d_{2}\right)^{2} \dot{\theta}_{2}^{2}+\frac{1}{2} I_{2} \dot{\theta}_{2}^{2}+m_{2} l_{1}\left(l_{2}-d_{2}\right) \dot{\theta}_{1} \dot{\theta}_{2} \cos \left(\theta_{1}-\theta_{2}\right) \tag{A14}
\end{equation*}
$$

$$
\begin{align*}
K E_{3}=\frac{1}{2} m_{3} v_{c 3}^{2}= & \frac{1}{2} m_{3} l_{1}^{2} \dot{\theta}_{1}^{2}+\frac{1}{2} m_{3} l_{2}^{2} \dot{\theta}_{2}^{2}+\frac{1}{2} m_{3}\left(l_{3}-d_{3}\right)^{2} \dot{\theta}_{3}^{2}+\frac{1}{2} I_{3} \dot{\theta}_{3}^{2} \\
& +m_{3} l_{1} l_{2} \dot{\theta}_{1} \dot{\theta}_{2} \cos \left(\theta_{1}-\theta_{2}\right)+m_{3} l_{1}\left(l_{3}-d_{3}\right) \dot{\theta}_{1} \dot{\theta}_{3} \cos \left(\theta_{1}-\theta_{3}\right)  \tag{A15}\\
& +m_{3} l_{2}\left(l_{3}-d_{3}\right) \dot{\theta}_{2} \dot{\theta}_{3} \cos \left(\theta_{2}-\theta_{3}\right)
\end{align*}
$$

The potential energy of each link:

$$
\begin{gather*}
P E_{1}=m_{1} g y_{1}=m_{1} g\left(l_{1}-d_{1}\right) \cos \theta_{1}  \tag{A16}\\
P E_{2}=m_{2} g y_{2}=m_{2} g l_{1} \cos \theta_{1}+m_{2} g\left(l_{2}-d_{2}\right) \cos \theta_{2}  \tag{A17}\\
P E_{3}=m_{3} g y_{3}=m_{3} g l_{1} \cos \theta_{1}+m_{3} g l_{2} \cos \theta_{2}+m_{3} g\left(l_{3}-d_{3}\right) \cos \theta_{3} \tag{A18}
\end{gather*}
$$

Now we derive the Lagrange equation of the model:

$$
\begin{equation*}
\mathcal{L}=\text { total kinetic energy - total potential energy }=K E-P E \tag{A19}
\end{equation*}
$$

$$
\begin{align*}
\mathcal{L}=\frac{1}{2}\left(m_{1}\left(l_{1}-d_{1}\right)^{2}\right. & \left.+m_{2} l_{1}^{2}+m_{3} l_{1}^{2}\right) \dot{\theta}_{1}^{2}+\frac{1}{2}\left(m_{2}\left(l_{2}-d_{2}\right)^{2}+m_{3} l_{2}^{2}\right) \dot{\theta}_{2}^{2}+\frac{1}{2} m_{3} d_{3}^{2} \dot{\theta}_{3}^{2} \\
& -\left(m_{2} l_{1}\left(l_{2}-d_{2}\right)+m_{3} l_{1} l_{2}\right) \dot{\theta}_{1} \dot{\theta}_{2} \cos \left(\theta_{1}+\theta_{2}\right) \\
& +m_{3} l_{1} d_{3} \dot{\theta}_{1} \dot{\theta}_{3} \cos \left(\theta_{1}-\theta_{3}\right)-m_{3} l_{2} d_{3} \dot{\theta}_{2} \dot{\theta}_{3} \cos \left(\theta_{2}+\theta_{3}\right)  \tag{A20}\\
& -g\left(m_{1}\left(l_{1}-d_{1}\right)+m_{2} l_{1}+m_{3} l_{1}\right) \cos \theta_{1} \\
& -g\left(m_{2}\left(l_{2}-d_{2}\right)+m_{3} l_{2}\right) \cos \theta_{2}-m_{3} g d_{3} \cos \theta_{3}
\end{align*}
$$

Now we convert the continuous Lagrange equation into a discrete Lagrange equation using the mid-point rule as follows:

$$
\begin{align*}
& \mathcal{L}(q, \dot{q}) \rightarrow h \mathcal{L}_{d}\left(\frac{q_{k+1}+q_{k}}{2}, \frac{q_{k+1}-q_{k}}{h}\right)  \tag{A21}\\
& \frac{1}{2 h} a_{2} \Delta \theta_{2 k}^{2}+\frac{1}{2 h} a_{3} \Delta \theta_{3 k}^{2}+\frac{1}{h} b_{1} \Delta \theta_{1 k} \Delta \theta_{2 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum}{2}\right.  \tag{A22}\\
& +\frac{1}{h} b_{2} \Delta \theta_{1 k} \Delta \theta_{3 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right) \\
& +\frac{1}{h} b_{3} \Delta \theta_{2 k} \Delta \theta_{3 k} \cos \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)-g h c_{1} \cos \left(\frac{\sum \theta_{1 k}}{2}\right)
\end{align*}
$$

$$
L_{d}=\frac{1}{2 h} a_{1} \Delta \theta_{1 k}^{2}+\frac{1}{2 h} a_{2} \Delta \theta_{2 k}^{2}+\frac{1}{2 h} a_{3} \Delta \theta_{3 k}^{2}+\frac{1}{h} b_{1} \Delta \theta_{1 k} \Delta \theta_{2 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right)
$$

$$
+\frac{1}{h} b_{2} \Delta \theta_{1 k} \Delta \theta_{3 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)
$$

$$
-g h c_{2} \cos \left(\frac{\sum \theta_{2 k}}{2}\right)-g h c_{3} \cos \left(\frac{\sum \theta_{3 k}}{2}\right)
$$

where:

$$
\begin{gather*}
a_{1}=I_{1}+m_{1}\left(l_{1}-d_{1}\right)^{2}+m_{2} l_{1}^{2}+m_{3} l_{1}^{2}  \tag{A23}\\
a_{2}=I_{2}+m_{2}\left(l_{2}-d_{2}\right)^{2}+m_{3} l_{2}^{2}  \tag{A24}\\
a_{3}=I_{3}+m_{3}\left(l_{3}-d_{3}\right)^{2}  \tag{A25}\\
b_{1}=m_{2} l_{1}\left(l_{2}-d_{2}\right)+m_{3} l_{1} l_{2}  \tag{A26}\\
b_{2}=m_{3} l_{1}\left(l_{3}-d_{3}\right)  \tag{A27}\\
b_{3}=m_{3} l_{2}\left(l_{3}-d_{3}\right)  \tag{A28}\\
c_{1}=m_{1}\left(l_{1}-d_{1}\right)+m_{2} l_{1}+m_{3} l_{1} \tag{A29}
\end{gather*}
$$

$$
\begin{equation*}
c_{2}=m_{2}\left(l_{2}-d_{2}\right)+m_{3} l_{2} \tag{A30}
\end{equation*}
$$

Now find $D_{1} L_{d}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)$ and $D_{2} L_{d}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)$. $D_{1} L_{d}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)$ : The first derivative of the discrete Lagrange with respect to the current coordinates (i.e., $\theta_{1 k}, \theta_{2 k}$, and $\theta_{3 k}$ )

$$
\begin{align*}
D_{1} \mathcal{L}_{d}=\frac{\partial \mathcal{L}_{d}}{\partial \theta_{1 k}} & \left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& =-\frac{1}{h} a_{1} \Delta \theta_{1 k}-\frac{1}{h} b_{1} \Delta \theta_{2 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right) \\
& -\frac{1}{2 h} b_{1} \Delta \theta_{1 k} \Delta \theta_{2 k} \sin \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right)-\frac{1}{h} b_{2} \Delta \theta_{3 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)  \tag{A31}\\
& -\frac{1}{2 h} b_{2} \Delta \theta_{1 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)+\frac{g h c_{1}}{2} \sin \left(\frac{\sum \theta_{1 k}}{2}\right) \\
D_{1} \mathcal{L}_{d}=\frac{\partial \mathcal{L}_{d}}{\partial \theta_{2 k}} \quad & \left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& =-\frac{1}{h} a_{2} \Delta \theta_{2 k}-\frac{1}{h} b_{1} \Delta \theta_{1 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right) \\
& +\frac{1}{2 h} b_{1} \Delta \theta_{1 k} \Delta \theta_{2 k} \sin \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right)-\frac{1}{h} b_{3} \Delta \theta_{3 k} \cos \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)  \tag{A32}\\
& -\frac{1}{2 h} b_{3} \Delta \theta_{2 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right) \frac{g h c_{2}}{2} \sin \left(\frac{\sum \theta_{2 k}}{2}\right) \\
D_{1} \mathcal{L}_{d}=\frac{\partial \mathcal{L}_{d}}{\partial \theta_{3 k}} & \left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& =-\frac{1}{h} a_{3} \Delta \theta_{3 k}-\frac{1}{h} b_{2} \Delta \theta_{1 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)  \tag{A33}\\
& +\frac{1}{2 h} b_{2} \Delta \theta_{1 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)-\frac{1}{h} b_{3} \Delta \theta_{2 k} \cos \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right) \\
& +\frac{1}{2 h} b_{3} \Delta \theta_{2 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)+\frac{g h c_{3}}{2} \sin \left(\frac{\sum \theta_{3 k}}{2}\right)
\end{align*}
$$

$D_{2} L_{d}\left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right)$ : The first derivative of the discrete Lagrange with respect to the current coordinates (i.e., $\theta_{1 k+1}, \theta_{2 k+1}$, and $\theta_{3 k+1}$ )

$$
\begin{align*}
D_{2} \mathcal{L}_{d}=\frac{\partial \mathcal{L}_{d}}{\partial \theta_{1 k+1}} & \left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& =\frac{1}{h} a_{1} \Delta \theta_{1 k}+\frac{1}{h} b_{1} \Delta \theta_{2 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right) \\
& -\frac{1}{2 h} b_{1} \Delta \theta_{1 k} \Delta \theta_{2 k} \sin \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right)+\frac{1}{h} b_{2} \Delta \theta_{3 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)  \tag{A34}\\
& -\frac{1}{2 h} b_{2} \Delta \theta_{1 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)+g h c_{1} \sin \left(\frac{\sum \theta_{1 k}}{2}\right) \\
D_{2} \mathcal{L}_{d}=\frac{\partial \mathcal{L}_{d}}{\partial \theta_{2 k+1}} \quad & \left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& =\frac{1}{h} a_{2} \Delta \theta_{2 k}+\frac{1}{h} b_{1} \Delta \theta_{1 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right) \\
& +\frac{1}{2 h} b_{1} \Delta \theta_{1 k} \Delta \theta_{2 k} \sin \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{2 k}}{2}\right)+\frac{1}{h} b_{3} \Delta \theta_{3 k} \cos \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)  \tag{A35}\\
& -\frac{1}{2 h} b_{3} \Delta \theta_{2 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)+g h c_{2} \sin \left(\frac{\sum \theta_{2 k}}{2}\right) \\
D_{2} \mathcal{L}_{d}=\frac{\partial \mathcal{L}_{d}}{\partial \theta_{3 k+1}} \quad & \left(\theta_{1 k}, \theta_{1 k+1}, \theta_{2 k}, \theta_{2 k+1}, \theta_{3 k}, \theta_{3 k+1}\right) \\
& =\frac{1}{h} a_{3} \Delta \theta_{3 k}+\frac{1}{h} b_{2} \Delta \theta_{1 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right) \\
& +\frac{1}{2 h} b_{2} \Delta \theta_{1 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)+\frac{1}{h} b_{3} \Delta \theta_{2 k} \cos \left(\frac{\sum \theta_{1 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)  \tag{A36}\\
& +\frac{1}{2 h} b_{3} \Delta \theta_{2 k} \Delta \theta_{3 k} \sin \left(\frac{\sum \theta_{2 k}}{2}-\frac{\sum \theta_{3 k}}{2}\right)+g h c_{3} \sin \left(\frac{\sum \theta_{3 k}}{2}\right)
\end{align*}
$$

After several investigations, we came out with the proper and most commonly used set of anthropometric information. Table A1 provides the weight of each segment as a percentage of the total body weight, the length of each segment as a percentage of total body height [25,26], the location of Center of Mass (COM) of each segment, measured as percentage of segment length [26], the segments' radius of gyration in the frontal plane (perpendicular to the sagittal plane) as a percentage of segment length [24], and it will be used to calculate the Moment of Inertia (MOI) for each segment. Based on this information, the required anthropometric data of the model is calculated for a male with 70 kg mass and 1.70 m height.

Table A1. Anthropometric percentages of different body segments required for the model.

| Body Segment | Segment Mass as a <br> PERCENTAGE of Body <br> Mass | Segment Length as a <br> Percentage of Total <br> Body Height | Distance of Segment <br> COM from Proximal <br> End as a Percentage of <br> Segment Length | Radius of Gyration of <br> Body Segments in Frontal <br> Plane as a Percentage of <br> Segment Length |
| :---: | :---: | :---: | :---: | :---: |
| Head and Neck | 8.2 | 10.75 | 56.7 | 31.5 |
| Trunk | 46.84 | 30.00 | 56.2 | 38.3 |
| Upper arm | 3.25 | 17.20 | 43.6 | 31.0 |
| Forearm | 1.8 | 15.70 | 43.0 | 28.4 |
| Hand | 0.65 | 5.75 | 46.8 | 23.3 |
| Thigh | 10.5 | 23.20 | 43.3 | 26.7 |
| Calf | 4.75 | 24.70 | 43.4 | 27.5 |

The moment of inertia for each segment is calculated using the anthropometric data associated with each segment. For instance, evaluation of MOI of the head and neck will be as follows:

- Mass of the head and neck segment is considered to be $8.2 \%$ of total body mass according to the data provided in Table A1, and hence, $m_{H N}=5.78 \mathrm{~kg}$.
- According to Table A1, the length of the head and neck segment is $10.75 \%$ of the body height and is found to be $L_{H N}=0.18 \mathrm{~m}$.
- Radius of gyration of the head and neck in the frontal axis is $31.5 \%$, and hence, $k_{H N}=0.058 \mathrm{~m}$.
- Now, the MOI of the head and neck can be determined using the following equation:

$$
\begin{equation*}
I_{H N}=m_{H N} k_{H N}^{2}=5.78 \times 0.058^{2}=0.0192 \mathrm{~kg} \cdot \mathrm{~m}^{2} \tag{A37}
\end{equation*}
$$

Using the same technique, the MOI of torso, hand segments, and leg segments are calculated, and the values are given in Table A2 with the other anthropometric measurements of each segment. The mass and moment of inertia of the head, the arms, and the torso (HAT) segment is simply found by summing up the individual mass of these segments; however, the center of mass is found using the parallel axis theorem [41]:

$$
\begin{equation*}
m_{H A T} d_{H A T}=m_{H N} d_{H N}+m_{T} d_{T}+2 * m_{U A} d_{U A}+2 * m_{F A} d_{F A}+2 * m_{H} d_{H} \tag{A38}
\end{equation*}
$$

The distance of the COM of the head and neck from the top of the head is $d_{H N}$ and is calculated directly from Table A1. The distance of the COM of the torso from the top of the head is $d_{T}$, and it is equivalent to the total length of the head and neck segment and the distance from the neck to the COM of the torso. The same procedure is followed to calculate the distance from the top of the head to the COM of each segment of the arm.

$$
\begin{align*}
(5.78+32.8+2 \times 2.275 & +2 \times 1.309+2 \times 0.455) d_{H A T} \\
& =5.78 \times 0.078+32.8 \times 0.4034+2 \times 2.275 \times 0.3075+2 \times 1.309 \times 0.585  \tag{A39}\\
& +2 \times 0.455 \times 0.84
\end{align*}
$$

$$
\begin{equation*}
d_{H A T}=\frac{17.379}{46.66}=0.3725 \mathrm{~m} \tag{A40}
\end{equation*}
$$



Figure A1. COM of body segments in the sagittal plane.
The MOI of upper body segments around the axis of rotation, which is the $y$-axis, can be approximated using the principle Parallel Axes Theorem. For applying the theorem, we need to determine the offset ( $r$ ) for each segment, which is the shortest distance between the axis of rotation and the other axis passing through the COG of that segment. For example, the offset distance between the axis of rotation and the axis of the head and neck segment is defined as $r_{H N}$, and it can be found as the distance between the COG of the head and neck segment to the joint of the shoulder and the distance from the joint of the shoulder to the COG of the HAT segment.

$$
\begin{equation*}
r_{H N}=d_{H A T}-d_{H N}=0.3725-0.078=0.2945 \mathrm{~m} \tag{A41}
\end{equation*}
$$

The MOI of the head and neck about the axis of rotation is found as shown below.

$$
\begin{equation*}
I_{H N}^{\prime}=I_{H N}+m_{H N} r_{H N}^{2}=0.0192+\left(5.78 \times 0.2945^{2}\right)=0.5205 \mathrm{~kg} \cdot \mathrm{~m}^{2} \tag{A42}
\end{equation*}
$$

Using the same procedure, we found the offset distances and the moments of inertia about the axis of rotation for the other segments:

For the torso:

$$
\begin{gather*}
r_{T}=d_{T}-d_{H A T}=0.4034-0.3725=0.0309 \mathrm{~m}  \tag{A43}\\
I_{T}^{\prime}=I_{T}+m_{T} r_{T}^{2}=1.251+\left(32.8 \times 0.0309^{2}\right)=1.2823 \mathrm{~kg} \cdot \mathrm{~m}^{2} \tag{A44}
\end{gather*}
$$

For the upper arms:

$$
\begin{equation*}
r_{U A}=d_{H A T}-d_{U A}=0.3725-0.3075=0.065 \mathrm{~m} \tag{A45}
\end{equation*}
$$

$$
\begin{equation*}
I_{U A}{ }^{\prime}=I_{U A}+m_{U A} r_{U A}^{2}=0.0187+\left(2.275 \times 0.065^{2}\right)=0.0283 \mathrm{~kg} \cdot \mathrm{~m}^{2} \tag{A46}
\end{equation*}
$$

For the forearms:

$$
\begin{gather*}
r_{F A}=d_{F A}-d_{H A T}=0.585-0.3075=0.2775 \mathrm{~m}  \tag{A47}\\
I_{F A}^{\prime}=I_{F A}+m_{F A} r_{F A}^{2}=0.00752+\left(1.309 \times 0.2775^{2}\right)=0.1083 \mathrm{~kg} \cdot \mathrm{~m}^{2} \tag{A48}
\end{gather*}
$$

For the hands:

$$
\begin{gather*}
r_{H}=d_{H}-d_{H A T}=0.84-0.3075=0.5325 \mathrm{~m}  \tag{A49}\\
I_{H}^{\prime}=I_{H}+m_{H} r_{H}^{2}=0.000236+\left(0.455 \times 0.5325^{2}\right)=0.1292 \mathrm{~kg} \cdot \mathrm{~m}^{2} \tag{A50}
\end{gather*}
$$

Then, the total MOI of the HAT segment about the axis of rotation is found by summing up the MOI of the individual segments about the same axis of rotation:

$$
\begin{gather*}
I_{H A T}{ }^{\prime}=I_{H N}{ }^{\prime}+I_{T}{ }^{\prime}+I_{U A}{ }^{\prime}+I_{F A}{ }^{\prime}+I_{H}{ }^{\prime} \\
=0.5205+1.2823+0.0283+0.1083+0.1292=2.0686 \mathrm{~kg} \cdot \mathrm{~m}^{2} \tag{A51}
\end{gather*}
$$

The table below shows the physical parameters for a male with 70 kg and 1.70 m required to build the model. Since the model does not include the head and the arms, the head and the three segments of the arms will be considered as one segment with the torso, represented as HAT (Head, Arms, and Torso).

Table A2. Model physical parameters based on a male with 70 kg mass and 1.70 m height.

| Body Segment | Segment Mass <br> $(\mathbf{K G})$ | Segment Length <br> $(\mathbf{m})$ | Distance of Segment <br> COM from Proximal <br> $(\mathbf{m})$ | Moment of Inertia |
| :---: | :---: | :---: | :---: | :---: |
| Head \& Neck | 0.18 | 5.782 | 0.10361925 | 0.019160814 |
| Trunk | 0.51 | 32.788 | 0.28662 | 1.250987086 |
| Upper arm | 0.29 | 2.275 | 0.1274864 | 0.018692162 |
| Forearm | 0.27 | 1.309 | 0.1155677 | 0.007520963 |
| Hand | 0.10 | 0.455 | 0.045747 | 0.000236024 |
| Thigh | 0.39 | 7.35 | 0.1707752 | 0.081504892 |
| Calf | 0.42 | 3.325 | 0.1822366 | 0.044335212 |
| HAT | 0.69 | 46.648 | 0.214937 | 2.0686 |

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# Efficient Method to Solve the Monge-Kantarovich Problem Using Wavelet Analysis 

Juan Rafael Acosta-Portilla ${ }^{1(D)}$, Carlos González-Flores ${ }^{2}{ }^{(\mathbb{D}}$, Raquiel Rufino López-Martínez ${ }^{3}$ (D) and Armando Sánchez-Nungaray $3, *$ (1)<br>1 Instituto de Investigaciones y Estudios Superiores Económicos y Sociales, Universidad Veracruzana, Veracruz 94294, Mexico; juaacosta@uv.mx<br>2 Escuela Superior de Ingeniería Mecánica y Eléctrica Zacatenco, Instituto Politécnico Nacional, Mexico City 07738, Mexico; cfgonzalez@esimez.mx<br>3 Facultad de Matemáticas, Universidad Veracruzana, Veracruz 94294, Mexico; ralopez@uv.mx<br>* Correspondence: armsanchez@uv.mx

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#### Abstract

In this paper, we present and justify a methodology to solve the Monge-Kantorovich mass transfer problem through Haar multiresolution analysis and wavelet transform with the advantage of requiring a reduced number of operations to carry out. The methodology has the following steps. We apply wavelet analysis on a discretization of the cost function level $j$ and obtain four components comprising one corresponding to a low-pass filter plus three from a high-pass filter. We obtain the solution corresponding to the low-pass component in level $j-1$ denoted by $\mu_{j-1}^{*}$, and using the information of the high-pass filter components, we get a solution in level $j$ denoted by $\hat{\mu}_{j}$. Finally, we make a local refinement of $\hat{\mu}_{j}$ and obtain the final solution $\mu_{j}^{\sigma}$.


Keywords: mass transfer problem; wavelets; multiresolution analysis

MSC: 42C40; 49Q20; 65T60

## 1. Introduction

In recent years, schemes to approximate infinite linear programs have become very important in theory. The authors of [1] showed that under suitable assumptions, the program's optimum value can be approximated by the values of finite-dimensional linear programs and that every accumulation point of a sequence of optimal solutions for the approximating programs is an optimal solution for the original problem. In particular, in [2] the authors studied the Monge-Kantorovich mass transfer (MT) problem on metric spaces. They considered conditions under the MT problem as solvable and, furthermore, that an optimal solution can be obtained as the weak limit of a sequence of optimal solutions to suitably approximate MT problems.

Moreover, in [3], the authors presented a numerical approximation for the value of the mass transfer (MT) problem on compact metric spaces. A sequence of transportation problems was built, and it proved that the value of the MT problem is a limit of the optimal values of these problems. Moreover, they gave an error bound for the numerical approximation. A generalization of this scheme of approximation was presented in $[4,5]$. They proposed an approximation scheme for the Monge-Kantorovich (MK) mass transfer problem on compact spaces that consisted of reducing to solve a sequence of finite transport problems. The method presented in that work uses a metaheuristic algorithm inspired by a scatter search in order to reduce the dimensionality of each transport problem. Finally, they provided some examples of that method.

On the other hand, the authors of [6] provided orthonormal bases for $L^{2}\left(\mathbb{R}^{n}\right)$ that have properties that are similar to those enjoyed by the classical Haar basis for $L^{2}(\mathbb{R})$. For example, each basis consists of appropriate dilates and translates of a finite collection of
"piecewise constant" functions. The construction is based on the notion of multiresolution analysis and reveals an interesting connection between the theory of compactly supported wavelet bases and the theory of self-similar tilings. Recent applications of the wavelet filter methodology have been used in various problems arising in communication systems and detection of thermal defects (see, for example, $[7,8]$, respectively).

In [9], the authors gave a scheme to approximate the MK problem based on the symmetries of the underlying spaces. They took a Haar-type MRA constructed according to the geometry of the spaces. Thus, they applied the Haar-type MRA based on symmetries to the MK problem and obtained a sequence-of-transport problem that approximates the original MK problem for each MRA space. Note that in the case of Haar's classical wavelet, this methodology coincides with the methods presented in [2,3].

It is important to note that various scientific problems are modeled through the Monge-Kantorovich approach; therefore, providing new efficient methodologies to find approximations of such problems turns out to be very useful. Within the applications of problems whose solutions are the Monge-Kantorovich problem are found: the use of the transport problem for the analysis of elastic image registration (see, for example, [10-12]). Other optimization problems related to this topic and differential equation tools can be found in recent works such as $[13,14]$.

The main goal of this paper is to present a scheme of approximation of the MK problem based on wavelet analysis in which we use wavelet filters to split the original problem. That is, we apply the filter to the discrete cost function in level $j$, which results in a cost function of level $j-1$ and three components of wavelet analysis. Using the information of the cost function given by the low-pass filter, which belongs to level $j-1$, we construct $\mu_{j-1}^{*}$ a solution of the MK problem for that level $j-1$, and using the additional information, the other three components of wavelet analysis are extended to $\hat{\mu}_{j}$, which is a solution to level $j$, where the projection of $\hat{\mu}_{j}$ to level $j-1$ is $\mu_{j-1}^{*}$. Finally, we make a local analysis of the solution $\hat{\mu}_{j}$ to obtain an improved solution based on the type of points of that solution (we have two type of points that are defined in the base in the connectedness of the solution).

This work has three non-introductory sections. In the first of them we present the Haar multiresolution analysis (MRA) in one and two dimensions. Next, we relate this to the absolutely continuous measures over a compact in $\mathbb{R}^{2}$. We finish with the definition of the Monge-Kantorovich mass transfer problem and its relation to the MRA.

In the second section, we define a proximity criterion for the components of the support of the simple solutions of the MK problem and study in detail the problem of, given a solution $\mu_{j-1}^{*}$ at level $j-1$ of resolution for the MK problem, construct a feasible solution $\hat{\mu}_{j}$ for the MK problem at level of resolution $j$ such that it is a refinement of the solution with lees resolution.

On the other hand, in the third section we present a methodological proposal to solve the $\mathrm{MK}^{j}$ problem such that it can be summarized in a simple algorithm of six steps:
Step 1. We consider a discretization of the cost function for the level $j$, denoted by $c_{j}$.
Step 2. We apply the wavelet transform to $c_{j}$; we obtain the low-pass component $c_{j-1}$ and three high-pass components, denoted by $\Psi_{1}, \Psi_{2}$ and $\Psi_{3}$, respectively.
Step 3. Using $c_{j-1}$ and the methodology of [3,4,9], we obtain a solution $\mu_{j-1}^{*}$ for $M K^{j-1}$ associated with this cost function.
Step 4. We classify the points of the support of the solution $\mu_{j-1}^{*}$ by proximity criteria as points of Type I or Type II.
Step 5. Using the solution $\mu_{j-1}^{*}$, the information of the high-pass components and Lemma 1, we obtain a feasible solution for the level $j$, which is denoted by $\hat{\mu}_{j}$. This feasible solution has the property that its projection to the level $j-1$ is equal to $\mu_{j-1}^{*}$; moreover, the support of $\hat{\mu}_{j}$ is contained in the support of $\mu_{j-1}^{*}$.
Step 6. The classification of the points of $\mu_{j-1}^{*}$ induce classification of the points in $\hat{\mu}_{j}$ by contention in the support. Over the points of Type I of the solution $\hat{\mu}_{j}$, we do not move those points. For the points of Type II, we apply a permutation to the
solution over the two points that better improves the solution, and we repeat the process with the rest of the points.
Finally, we present a series of examples that use the proposed methodology based on wavelet analysis and compare their results with those obtained applying the methodology of $[3,4,9]$.

## 2. Preliminaries

### 2.1. One-Dimensional MRA

The results of this and the following subsection are well known, and for a detailed exposition, we recommend consulting [15-17]. We begin by defining a general multiresolution analysis and developing the particular case of the Haar multiresolution analysis on $\mathbb{R}$. Given $a>0$ and $b \in \mathbb{R}$, the dilatation operator $D_{a}$ and the translation operator $T_{b}$ are defined by

$$
\begin{equation*}
\left(D_{a} f\right)(x)=a^{1 / 2} f(a x) \quad \text { and } \quad\left(T_{b} f\right)(x)=f(x-b) \tag{1}
\end{equation*}
$$

for every $f \in L^{2}(\mathbb{R})$, where the latter denotes the usual Hilbert space of square integrable real functions defined on $\mathbb{R}$. A multiresolution analysis (MRA) on $\mathbb{R}$ is a sequence of subspaces $\left(V_{j}\right)_{j \in \mathbb{Z}}$ of $L^{2}(\mathbb{R})$ such that it satisfies the following properties:
(1) $V_{j} \subset V_{j+1}$ for every $j \in \mathbb{Z}$.
(2) $\quad L^{2}(\mathbb{R})=\overline{\operatorname{span}} \bigcup_{j \in \mathbb{Z}} V_{j}$.
(3) $\bigcap_{j \in \mathbb{Z}} V_{j}=\{0\}$.
(4) $V_{j}=D_{2^{j}} V_{0}$.
(5) There exists a function $\varphi \in L^{2}$, called the scaling function, such that the collection $\left\{T_{j} \varphi\right\}_{j \in \mathbb{Z}}$ is an orthonormal system of translates and

$$
\begin{equation*}
V_{0}=\overline{\operatorname{span}}\left\{T_{j} \varphi\right\}_{i \in \mathbb{Z}} . \tag{2}
\end{equation*}
$$

We denote as $\chi_{A}$ the characteristic function of the set $A$. Then, the Haar scaling function is defined by

$$
\begin{equation*}
\varphi(x)=\chi_{[0,1)}(x) \tag{3}
\end{equation*}
$$

For each pair $j, k \in \mathbb{Z}$, we call

$$
\begin{equation*}
I_{j, k}=\left[2^{-j} k, 2^{-j}(k+1)\right) . \tag{4}
\end{equation*}
$$

Hence, we define the function

$$
\begin{align*}
\varphi_{j, k}(x) & =\left(T_{k} D_{2^{j}} \varphi\right)(x) \\
& =2^{j / 2} \varphi\left(2^{j} x-k\right)  \tag{5}\\
& =2^{j / 2} \chi_{I_{j, k}}(x) .
\end{align*}
$$

The collection $\left\{\varphi_{j, k}\right\}_{j, k \in \mathbb{Z}}$ is called the system of Haar scaling functions. For $j_{0} \in \mathbb{Z}$, the collection $\left\{\varphi_{j_{0}, k}\right\}_{k \in \mathbb{Z}}$ is referred to as the system of scale $j_{0}$ Haar scaling functions. The Haar function is defined by

$$
\begin{equation*}
\psi(x)=\chi_{[0,1 / 2)}(x)-\chi_{[1 / 2,1)}(x) \tag{6}
\end{equation*}
$$

For each pair $j, k \in \mathbb{Z}$, we define the function

$$
\begin{align*}
\psi_{j, k}(x) & =\left(T_{k} D_{2^{j}} \psi\right)(x) \\
& =2^{j / 2} \psi\left(2^{j} x-k\right)  \tag{7}\\
& =2^{j / 2}\left(\chi_{I_{j+1,2 k}}-\chi_{I_{j+1,2 k+1}}\right)(x)
\end{align*}
$$

The collection $\left\{\psi_{j, k}\right\}_{j, k \in \mathbb{Z}}$ is referred to as the Haar system on $\mathbb{R}$. For $j_{0} \in \mathbb{Z}$, the collection $\left\{\psi_{j_{0}, k}\right\}_{k \in \mathbb{Z}}$ is referred to as the system of scale $j_{0}$ Haar functions. It is well known that with respect to the usual inner product $\langle\cdot, \cdot\rangle$ in $L^{2}(\mathbb{R})$, the Haar system on $\mathbb{R}$ is an orthonormal system. Moreover, for each $j_{0} \in \mathbb{Z}$, the collection of scale $j_{0}$ Haar scaling functions is an orthonormal system. Thus, for each $j \in \mathbb{Z}$, the approximation operator $P_{j}$ on $L^{2}(\mathbb{R})$ is defined by

$$
\begin{equation*}
\left(P_{j} f\right)(x)=\sum_{k}\left\langle f, \varphi_{j, k}\right\rangle \varphi_{j, k}(x), \text { for all } f \in L^{2}(\mathbb{R}) \tag{8}
\end{equation*}
$$

and the approximation space $V_{j}$ by

$$
\begin{equation*}
V_{j}=\overline{\operatorname{span}}\left\{\varphi_{j, k}\right\}_{k \in \mathbb{Z}} . \tag{9}
\end{equation*}
$$

The collection $\left\{V_{j}\right\}_{j \in \mathbb{Z}}$ is called Haar multiresolution analysis. Similarly, we have that for each $j_{0} \in \mathbb{Z}$, the detail operator $Q_{j}$ on $L^{2}(\mathbb{R})$ is defined by

$$
\begin{equation*}
\left(Q_{j} f\right)(x)=\left(P_{j+1} f\right)(x)-\left(P_{j} f\right)(x), \text { for all } f \in L^{2}(\mathbb{R}) \tag{10}
\end{equation*}
$$

and the wavelet space $W_{j}$ by

$$
\begin{equation*}
W_{j}=\overline{\operatorname{span}}\left\{\psi_{j, k}\right\}_{k \in \mathbb{Z}} . \tag{11}
\end{equation*}
$$

Note that $V_{j+1}=V_{j} \oplus W_{j}$ for all $j \in \mathbb{Z}$. Hence, the collection

$$
\begin{equation*}
\left\{\varphi_{j_{0}, k}, \psi_{j, k}: j \geq j_{0}, k \in \mathbb{Z}\right\} \tag{12}
\end{equation*}
$$

is a complete orthonormal system on $\mathbb{R}$; this system is called the scale $j_{0}$ Haar system on $\mathbb{R}$. As a consequence, the Haar system $\left\{\psi_{j, k}\right\}_{j, k \in \mathbb{Z}}$ is a complete orthonormal system on $\mathbb{R}$.

### 2.2. Two-Dimensional MRA

To obtain the Haar MRA on $\mathbb{R}^{2}$, we consider the Haar MRA on $\mathbb{R}$ defined in the previous subsection with scaling and Haar functions $\varphi$ and $\psi$, and from them, through a tensor product approach we can construct a two-dimensional scaling and Haar function. First, we define the four possible products:

$$
\begin{align*}
\Phi(x, y) & =\varphi(x) \varphi(y), \quad \Psi^{(1)}(x, y)=\varphi(x) \psi(y)  \tag{13}\\
\Psi^{(2)}(x, y) & =\psi(x) \varphi(y), \quad \Psi^{(3)}(x, y)=\psi(x) \psi(y)
\end{align*}
$$

which are the scaling function associated with the unitary square and three Haar functions, respectively. Hence, for each $j, k_{1}, k_{2} \in \mathbb{Z}$, we define naturally the scaling and Haar function systems:

$$
\begin{array}{ll}
\Phi_{j, k_{1}, k_{2}}(x, y)=\varphi_{j, k_{1}}(x) \varphi(y)_{j, k_{2}}, & \Psi_{j, k_{1}, k_{2}}^{(1)}(x, y)=\varphi(x)_{j, k_{1}} \psi(y)_{j, k_{2}} \\
\Psi_{j, k_{1}, k_{2}}^{(2)}(x, y)=\psi(x)_{j, k_{1}} \varphi(y)_{j, k_{2}}, & \Psi_{j, k_{1}, k_{2}}^{(3)}(x, y)=\psi_{j, k_{1}}(x) \psi(y)_{j, k_{2}} . \tag{14}
\end{array}
$$

Then for $j_{0} \in \mathbb{Z}$, as in the one-dimensional case, we have that the collection

$$
\begin{equation*}
\left\{\Phi_{j_{0}, k_{1}, k_{2}}: k_{1}, k_{2} \in \mathbb{Z}\right\} \cup\left\{\Psi_{j, k_{1}, k_{2}}^{(i)}: 1 \leq i \leq 3, j \geq j_{0}\right\} \tag{15}
\end{equation*}
$$

is an orthonormal basis on $L^{2}\left(\mathbb{R}^{2}\right)$. Thus, the collection

$$
\begin{equation*}
\left\{\Psi_{j, k_{1}, k_{2}}^{(i)}: 1 \leq i \leq 3, j, k_{1}, k_{2} \in \mathbb{Z}\right\} \tag{16}
\end{equation*}
$$

is an orthonormal basis on $L^{2}\left(\mathbb{R}^{2}\right)$. Then for $j \in \mathbb{Z}$ and $f \in L^{2}\left(\mathbb{R}^{2}\right)$, the approximation operator is defined by

$$
\begin{equation*}
\left(\mathbf{P}_{j} f\right)(x, y)=\sum_{k_{1}} \sum_{k_{2}}\left\langle f, \Phi_{j, k_{1}, k_{2}}\right\rangle \Phi_{j, k_{1}, k_{2}}(x, y) \tag{17}
\end{equation*}
$$

and for $i=1,2,3$, the detail operators are

$$
\begin{equation*}
\left(\mathbf{Q}_{j}^{(i)} f\right)(x, y)=\sum_{k_{1}} \sum_{k_{2}}\left\langle f, \Psi_{j, k_{1}, k_{2}}^{(i)}\right\rangle \Psi_{j, k_{1}, k_{2}}^{(i)}(x, y) . \tag{18}
\end{equation*}
$$

Hence, the projection can be written as

$$
\begin{equation*}
\left(\mathbf{P}_{j+1} f\right)(x, y)=\left(\mathbf{P}_{j} f\right)(x, y)+\left(\mathbf{Q}_{j}^{(1)} f\right)(x, y)+\left(\mathbf{Q}_{j}^{(2)} f\right)(x, y)+\left(\mathbf{Q}_{j}^{(3)} f\right)(x, y) \tag{19}
\end{equation*}
$$

We will describe the approximation $\mathbf{P}_{j}$ and detail operators $\mathbf{Q}_{j}^{(i)}$ from the geometric point of view. First of all, we fix some $j, k_{1}, k_{2} \in \mathbb{Z}$ and define the square

$$
\begin{align*}
S\left(j, k_{1}, k_{2}\right) & =I_{j, k_{1}} \times I_{j, k_{2}}  \tag{20}\\
& =\left[2^{-j} k_{1}, 2^{-j}\left(k_{1}+1\right)\right) \times\left[2^{-j} k_{2}, 2^{-j}\left(k_{2}+1\right)\right)
\end{align*}
$$

Then, we have that $\Phi_{j, k_{1}, k_{2}}$ is the characteristic function of $S\left(j, k_{1}, k_{2}\right)$, in symbols

$$
\begin{equation*}
\Phi_{j, k_{1}, k_{2}}(x, y)=\chi_{S}(x, y) \tag{21}
\end{equation*}
$$

where $S=S\left(j, k_{1}, k_{2}\right)$. Therefore for $f \in L^{2}\left(\mathbb{R}^{2}\right)$, the operator $\mathbf{P}_{j} f$ acts as a discretization of $f$ that is constant over the disjointed $S\left(j, k_{1}, k_{2}\right)$ squares. On the other hand, we can split $S\left(j, k_{1}, k_{2}\right)$ as follows (see Figure 1):


Figure 1. Functions $\Phi_{j, k_{1}, k_{2}}, \Psi_{j, k_{1}, k_{2}}^{(1)} \Psi_{j, k_{1}, k_{2}}^{(2)}$ and $\Psi_{j, k_{1}, k_{2}}^{(3)}$.
(i) Into two rectangles with half the height and the same width as $S\left(j, k_{1}, k_{2}\right)$, namely, the sets

$$
\begin{align*}
& A_{1}\left(j, k_{1}, k_{2}\right)=\left[2^{-j} k_{1}, 2^{-j}\left(k_{1}+1\right)\right) \times\left[2^{-j-1} 2 k_{2}, 2^{-j-1}\left(2 k_{2}+1\right)\right) \text { and } \\
& A_{2}\left(j, k_{1}, k_{2}\right)=\left[2^{-j} k_{1}, 2^{-j}\left(k_{1}+1\right)\right) \times\left[2^{-j-1} 2 k_{2}+1,2^{-j-1}\left(2 k_{2}+2\right)\right) . \tag{22}
\end{align*}
$$

(ii) Into two rectangles with the same height and half the width as $S\left(j, k_{1}, k_{2}\right)$, namely,

$$
\begin{align*}
& B_{1}\left(j, k_{1}, k_{2}\right)=\left[2^{-j-1} 2 k_{1}, 2^{-j-1}\left(2 k_{1}+1\right)\right) \times\left[2^{-j} k_{2}, 2^{-j}\left(k_{2}+1\right)\right) \text { and } \\
& B_{2}\left(j, k_{1}, k_{2}\right)=\left[2^{-j-1} 2 k_{1}+1,2^{-j-1}\left(2 k_{1}+2\right)\right) \times\left[2^{-j} k_{2}, 2^{-j}\left(k_{2}+1\right)\right) . \tag{23}
\end{align*}
$$

(iii) Into four squares with half the side lengths of $S\left(j, k_{1}, k_{2}\right)$, namely,

$$
\begin{align*}
& C_{1}\left(j, k_{1}, k_{2}\right)=\left[2^{-j-1} 2 k_{1}, 2^{-j-1}\left(2 k_{1}+1\right)\right) \times\left[2^{-j-1} 2 k_{2}, 2^{-j-1}\left(2 k_{2}+1\right)\right), \\
& C_{2}\left(j, k_{1}, k_{2}\right)=\left[2^{-j-1} 2 k_{1}+1,2^{-j-1}\left(2 k_{1}+2\right)\right) \times\left[2^{-j-1} 2 k_{2}, 2^{-j-1}\left(2 k_{2}+1\right)\right),  \tag{24}\\
& C_{3}\left(j, k_{1}, k_{2}\right)=\left[2^{-j-1} 2 k_{1}, 2^{-j-1}\left(2 k_{1}+1\right)\right) \times\left[2^{-j-1} 2 k_{2}+1,2^{-j-1}\left(2 k_{2}+2\right)\right) \text { and } \\
& C_{4}\left(j, k_{1}, k_{2}\right)=\left[2^{-j-1} 2 k_{1}+1,2^{-j-1}\left(2 k_{1}+2\right)\right) \times\left[2^{-j-1} 2 k_{2}+1,2^{-j-1}\left(2 k_{2}+2\right)\right) .
\end{align*}
$$

Hence, for $i=1,2,3$, the function $\Psi_{j, k_{1}, k_{2}}^{(i)}$ is defined by

$$
\begin{align*}
& \Psi_{j, k_{1}, k_{2}}^{(1)}(x, y)=\chi_{A_{1}}(x, y)-\chi_{A_{2}}(x, y) \text { with } A_{l}=A_{l}\left(j, k_{1}, k_{2}\right) \text { for } l=1,2 . \\
& \Psi_{j, k_{1}, k_{2}}^{(2)}(x, y)=\chi_{B_{1}}(x, y)-\chi_{B_{2}}(x, y) \text { with } B_{l}=B_{l}\left(j, k_{1}, k_{2}\right) \text { for } l=1,2 .  \tag{25}\\
& \Psi_{j, k_{1}, k_{2}}^{(3)}(x, y)=\chi_{C_{1} \cup C_{4}}(x, y)-\chi_{C_{2} \cup C_{3}}(x, y) \text { with } C_{l}=C_{l}\left(j, k_{1}, k_{2}\right) \text { for } l=1,2,3,4 .
\end{align*}
$$

(see Figure 1). Thus, the image of $f \in L^{2}\left(\mathbb{R}^{2}\right)$ under the detail operator $\mathbf{Q}_{j}^{(i)}$ is a function $\mathbf{Q}_{j}^{(i)} f$ formed by pieces that oscillate symmetrically on each square $S\left(j, k_{1}, k_{2}\right)$.

Now we have the elements to define an MRA on $\mathbb{R}^{2}$; to do this, we will use the notation introduced for the one-dimensional MRA on $\mathbb{R}$ defined in the previous subsection. For each $j \in \mathbb{Z}$, we define

$$
\begin{equation*}
\mathbf{V}_{j}=V_{j} \otimes V_{j} \tag{26}
\end{equation*}
$$

Then the collection $\left\{\mathbf{V}_{j}\right\}_{j \in \mathbb{Z}}$ is the Haar multiresolution analysis on $\mathbb{R}^{2}$, where the dilatation and translation operators are defined by

$$
\begin{equation*}
\mathbf{D}_{2 j}=D_{2^{j}} \otimes D_{2^{j}} \text { and } \mathbf{T}_{m, n}=T_{m} \otimes T_{n} \tag{27}
\end{equation*}
$$

respectively. Note that we have the following relation:

$$
\begin{equation*}
\mathbf{V}_{j+1}=\mathbf{V}_{j} \oplus \mathbf{W}_{j}^{(1)} \oplus \mathbf{W}_{j}^{(2)} \oplus \mathbf{W}_{j}^{(3)} \tag{28}
\end{equation*}
$$

where $\mathbf{W}_{j}^{(1)}=\mathbf{V}_{j} \otimes \mathbf{W}_{j}, \mathbf{W}_{j}^{(2)}=\mathbf{W}_{j} \otimes \mathbf{V}_{j}$ and $\mathbf{W}_{j}^{(3)}=\mathbf{W}_{j} \otimes \mathbf{W}_{j}$. For more detail with respect to the Haar MRA, see [15].

### 2.3. Measures and MRA

In this subsection, we will use the two previous ones and the [9] approach to relate measures over $\mathbb{R}^{2}$ with the Haar MRA on $\mathbb{R}^{2}$. The results and definitions presented in this and the following subsection can be found in the classical references [18,19]. We consider a compact subset $X$ of $\mathbb{R}^{2}$ and a measure $\mu$ such that it is absolutely continuous with respect to the Lebesgue measure $\lambda$. We call

$$
\begin{equation*}
f_{\mu}=\frac{d \mu}{d \lambda} \tag{29}
\end{equation*}
$$

the Radon-Nikodym derivative of $\mu$ with respect to $\lambda$. By construction, it necessary holds that $f_{\mu} \in L^{1}(X)$. We additionally suppose that $f_{\mu} \in L^{2}(X)$. Then, as a consequence of the Haar MRA on $\mathbb{R}^{2}$, we have that

$$
\begin{equation*}
\left\|\mathbf{P}_{j} f_{\mu}-f_{\mu}\right\|_{2} \rightarrow 0 \text { as } j \rightarrow \infty . \tag{30}
\end{equation*}
$$

Moreover, the compactness of $X$ ensures that

$$
\begin{equation*}
\left\|\mathbf{P}_{j} f_{\mu}-f_{\mu}\right\|_{1} \leq \lambda(X)^{1 / 2}\left\|\mathbf{P}_{j} f_{\mu}-f_{\mu}\right\|_{2} \tag{31}
\end{equation*}
$$

The above allow us to define $\mu_{j}$, the approximation of the measure $\mu$ to the level $j \in \mathbb{Z}$ of the Haar MRA on $\mathbb{R}^{2}$, as the measure induced by the projection of the Radon-Nikodym to the level $j$. That is, $\mu_{j}$ is defined by the relation

$$
\begin{equation*}
d \mu_{j}=\mathbf{P}_{j} f_{\mu} d \lambda \tag{32}
\end{equation*}
$$

We denote the expectation of a function $g \in L\left(\mathbb{R}^{2}\right)$ with respect to $\lambda$ as

$$
\begin{equation*}
\mathbb{E}[g]=\int_{\mathbb{R}^{2}} g d \lambda \tag{33}
\end{equation*}
$$

Then Theorem 4 and Corollary 5 in [9] ensure that for each $g \in L^{2}\left(\mathbb{R}^{2}\right)$ and $j \in \mathbb{Z}$, it is fulfilled that

$$
\begin{equation*}
\mathbb{E}\left[\mathbf{P}_{j} g\right]=\mathbb{E}[g] \tag{34}
\end{equation*}
$$

Thus, $\mu_{j}$ is absolutely continuous with respect to the Lebesgue measure. If, additionally, we suppose that $\mu$ is a measure with support on $X$, then $\left(\mu_{j}\right)$ converges to $\mu$ in the $L^{1}$ and $L^{2}$ sense. That is, by the Riesz Theorem, we can associate each of these measures $\mu_{j}$ to an integrable function $\hat{\mu}_{j} \in \mathbf{P}_{j} L^{2}\left(\mathbb{R}^{2}\right)$ such that for each Lebesgue measurable set $A \subset \mathbb{R}^{2}$,

$$
\begin{align*}
\mu_{j}(A) & =\int_{A} \hat{\mu}_{j} d \lambda \\
& =\int \hat{\mu}_{j} \chi_{A} d \lambda  \tag{35}\\
& =\mathbb{E}\left[\hat{\mu}_{j} \chi_{A}\right]
\end{align*}
$$

and we apply the respective convergence mode. Further, the compact support of measures ensures that the sequence $\left(\mu_{j}\right)$ converges weakly to the measure $\mu$, proof of which can be found in [9] Theorem 7.

### 2.4. M-K Problem and MRA

In this subsection, we study the Monge-Kantorovich problem from the point of view of the Haar MRA on $\mathbb{R}^{2}$ (for a detailed exposition, we recommend consulting [9]). Let $X$ and $Y$ be two compact subsets of $\mathbb{R}^{2}$. We denote by $M^{+}(X \times Y)$ the family of finite measures on $X \times Y$. Given $\mu \in M^{+}(X \times Y)$, we denote its marginal measures on $X$ and $Y$ as

$$
\begin{equation*}
\Pi_{1} \mu\left(E_{1}\right)=\mu\left(E_{1} \times Y\right) \tag{36}
\end{equation*}
$$

and

$$
\begin{equation*}
\Pi_{2} \mu\left(E_{2}\right)=\mu\left(X \times E_{2}\right) \tag{37}
\end{equation*}
$$

for each $\mu$-measurable set $E_{1} \subset X$ and $E_{2} \subset Y$. Let $c$ be a real function defined on $X \times Y$, and $\eta^{1}, \eta^{2}$ are two measures defined on $X$ and $Y$, respectively. The Monge-Kantorovich mass transfer problem is given as follows:

$$
\begin{align*}
\text { MK: minimize } & \langle\mu, c\rangle:=\int c d \mu  \tag{38}\\
\text { subject to: } & \Pi_{1} \mu=\eta^{1}, \quad \Pi_{2} \mu=\eta^{2}, \quad \mu \in M^{+}(X \times Y) .
\end{align*}
$$

A measure $\mu \in M^{+}(X \times Y)$ is said to be a feasible solution for the MK problem if it satisfies (38) and $\langle\mu, c\rangle$ is finite. We said that the MK problem is solvable if there is a feasible solution $\mu^{*}$ that attains the optimal value for it. So $\mu^{*}$ is called an optimal solution for (38). If, additionally, we assume that $\mu, \eta^{1}$ and $\eta^{2}$ are absolutely continuous with respect to the Lebesgue measures on $\mathbb{R}^{2}$ and $\mathbb{R}$, then in a natural way, we can discretize the MK problem
through the Haar MRA on $\mathbb{R}^{2}$ as follows. For $j \in \mathbb{Z}$, we define the MK problem of level $j$ by:

$$
\begin{align*}
\mathrm{MK}_{j}: \text { minimize } & \left\langle\mu_{j}, c\right\rangle:=\int c d \mu_{j}  \tag{39}\\
\text { subject to: } & \Pi_{1} \mu_{j}=\eta_{j}^{1}, \quad \Pi_{2} \mu_{j}=\eta_{j}^{2}, \quad \mu \in M^{+}(X \times Y)
\end{align*}
$$

where $\mu_{j}, \eta_{j}^{1}$ and $\eta_{j}^{2}$ are the projections to level $j$ of the measures $\mu, \eta^{1}$ and $\eta^{2}$, respectively, to the Haar MRA.

## 3. Technical Results

In this section, we present a series of results that ensure the good behavior of the methodological proposal of the next section. In order to do this, we start by assuming an MK problem with cost function $c=c(x, y)$, base sets $X=Y=[0,1]$ and measure restrictions $\eta^{1}=\eta^{2}=\left.\lambda\right|_{[0,1]}$. In other words, we consider the problem of moving a uniform distribution to a uniform one with the minimum movement cost. Since in applications we work with discretized problems, then as a result of applying the MRA on $\mathbb{R}^{2}$, we have that our objective is to solve:

$$
\begin{align*}
\mathrm{MK}_{j}: \text { minimize } & \sum_{i, k} \mu_{i, k}^{j} c_{i, k}^{j} \\
\text { subject to: } & \sum_{i} \mu_{i, k}^{j}=\frac{1}{2^{j}}, \text { for each } \quad k=1, \ldots, 2^{j} . \\
& \sum_{k} \mu_{i, k}^{j}=\frac{1}{2^{j}}, \text { for each } \quad i=1, \ldots, 2^{j} .  \tag{40}\\
& \sum_{i, k} \mu_{i, k}^{j}=1 . \\
& \mu_{i, k}^{j} \geq 0, \text { for each } i, k=1, \ldots, 2^{j} .
\end{align*}
$$

where $\mu_{i, k}^{j}$ is the portion of the initial value $\frac{1}{2^{j}}$ in the position $I_{j, i}$ of the $x$-axis allocated to the position $I_{j, k}$ of the $y$-axis. We call $j$-discrete unit square the grid formed by the squares $S\left(j, k_{1}, k_{2}\right)$ (see (20)), dividing the set $[0,1] \times[0,1]$ in $2^{j} \times 2^{j}$ blocks, in a such way that each one is identified with the point $\left(k_{1}, k_{2}\right)$. We suppose that there is a simple solution $\mu$ for (40). That is, $\mu$ is a feasible solution such that, given $i_{0}, k_{0} \in\left\{1, \ldots, 2^{j}\right\}$ with $\mu_{i_{0}, k_{0}} \neq 0$, it necessarily holds that $\mu_{i_{0}, k}=0$ for each $k \neq k_{0}$ and $\mu_{i, k_{0}}=0$ for each $i \neq i_{0}$. Geometrically, if the measure $\mu$ is plotted as a discrete heat map in the $j$-discrete unit square, then no color element in the plot has another color element in its same row and column, as can be seen in Figure 2.

Definition 1. We define a proximity criteria in the $j$-discrete unit square as follows: $(i, k)$ is a neighbor of $(1, m)$ if

$$
\begin{equation*}
|i-l|=1 \text { or }|k-m|=1 \tag{41}
\end{equation*}
$$

In Figure 2, we plot the support of the hypothetical simple solution $\mu$. Hence, the neighbors of the position in the middle of the cross are those that touch the yellow stripes. Then in this example, the middle point has four neighbors.


Figure 2. Support of $\mu$ and the proximity criteria.
With this in mind, we can classify the points in $\operatorname{support}(\mu)$ as follows.

Definition 2. We say that $\left(k_{1}, k_{2}\right) \in \operatorname{support}(\mu)$ is a border point if $k_{1}$ or $k_{2}$ equals 0 or 1 ; otherwise, we call it an interior point. It is clear that a border point has at least one neighbor and at most three, whereas an interior point has at least two neighbors and at most four. Hence, we can partition support ( $\mu$ ) into two sets as follows.

The set of the points of Type I is given by

$$
\begin{equation*}
T_{1}=\{(i, k) \in \operatorname{support}(\mu) \mid(i, k) \text { has minimun neighbors }\}, \tag{42}
\end{equation*}
$$

and the set of the points of Type II is given by

$$
\begin{equation*}
T_{2}=\operatorname{support}(\mu) \backslash T_{1} . \tag{43}
\end{equation*}
$$

Intuitively, the set $T_{1}$ is composed of well-controlled points, whereas the set $T_{2}$ has the points that admit permutations between them, since, as we will see in the next section, in the proposed algorithm they will be permuted. See Figures 3 and 4. Naturally, since $\mu$ is a feasible solution for (40), then given elements $\left(k_{1}, k_{2}\right),\left(m_{1}, m_{2}\right) \in \operatorname{support}(\mu)$ and for $i=1,2$ permutations $\sigma_{i}$ over $\left\{k_{i}, m_{i}\right\}$, the measure $\mu^{\sigma}$ defined by

$$
\begin{equation*}
\mu^{\sigma}(a, b)=\mu\left(\sigma_{1}(a), \sigma_{2}(b)\right) \tag{44}
\end{equation*}
$$

is a feasible solution.


Figure 3. Classification of the points in $\operatorname{support}(\mu)$ : Type I points.


Figure 4. Classification of the points in $\operatorname{support}(\mu)$ : Type II points.

## Refining Projections

In this subsection, we study the problem of improving an optimal solution $\mu_{j-1}^{*}$ for (40) on level $j-1 \in \mathbb{Z}$ to a feasible solution $\hat{\mu}_{j}$ for the next level $j$. Let $\mu_{j-1}^{*}$ be an optimal solution for level $j-1$. Then we are looking for $\hat{\mu}_{j}$ such that:
(1) $\hat{\mu}_{j}$ is a feasible solution.
(2) $\mathbf{P}_{j-1} \hat{\mu}_{j}=\mu_{j-1}^{*}$.

As described in the previous section, the measure $\hat{\mu}_{j} \in L^{2}\left(\mathbb{R}^{2}\right)$ can be decomposed in

$$
\begin{equation*}
\hat{\mu}_{j}=\mu_{j-1}^{*}+v_{1}+v_{2}+v_{3} \tag{46}
\end{equation*}
$$

where

$$
\begin{align*}
v_{i}(x, y) & =\left(\mathbf{Q}_{j}^{(i)} \hat{\mu}_{j}\right)(x, y) \\
& =\sum_{k_{1}} \sum_{k_{2}}\left\langle\hat{\mu}_{j}, \Psi_{j, k_{1}, k_{2}}^{(i)}\right\rangle \Psi_{j, k_{1}, k_{2}}^{(i)}(x, y) . \tag{47}
\end{align*}
$$

From the geometric point of view, the projections $v_{i}$ are formed from differences of characteristic functions, as we mentioned in Section 2.2. So we have the following result:

Lemma 1. Let $j \in \mathbb{Z}$ and $\mu_{j-1}^{*}$ be an optimal solution for (40) at level $j-1$. Then for each positive measure $\hat{\mu}_{j} \in \mathbf{P}_{j} L^{2}\left(\mathbb{R}^{2}\right)$ such that $\mathbf{P}_{j-1} \hat{\mu}_{j}=\mu_{j-1}^{*}$ and $\hat{\mu}_{j}=\mu_{j-1}^{*}+v_{1}+v_{2}+v_{3}$, it necessarily holds that

$$
\begin{equation*}
\mu_{j-1}^{*}(A)=0 \text { implies } v_{i}(A)=0 \tag{48}
\end{equation*}
$$

for each Lebesgue measurable set $A \subset \mathbb{R}^{2}$ and each $i=1,2,3$. Therefore, the support of $\hat{\mu}_{j}$ is contained in the support of $\mu_{j-1}^{*}$.

Proof. We only make the proof for the case $i=1$, since the other two are very similar. To simplify the notation, we use the symbols $\mu_{-}$and $v_{-}$as measures or functions in the respective subspace of $L^{2}\left(\mathbb{R}^{2}\right)$. Since when setting a level $j \in \mathbb{Z}$ all the measures in question are constant in pairs of rectangles dividing $S\left(j, k_{1}, k_{2}\right)$, as we prove in Section 2.2, then it is enough to prove that (65) is valid on this rectangles. Let

$$
\begin{align*}
A_{1} & =A_{1}\left(j, k_{1}^{\prime}, k_{2}^{\prime}\right) \\
& =\left[2^{-j} k_{1}^{\prime}, 2^{-j}\left(k_{1}^{\prime}+1\right)\right) \times\left[2^{-j-1} 2 k_{2}^{\prime}, 2^{-j-1}\left(2 k_{2}^{\prime}+1\right)\right) \tag{49}
\end{align*}
$$

and

$$
\begin{align*}
A_{2} & =A_{2}\left(j, k_{1}^{\prime}, k_{2}^{\prime}\right)  \tag{50}\\
& =\left[2^{-j} k_{1}^{\prime}, 2^{-j}\left(k_{1}^{\prime}+1\right)\right) \times\left[2^{-j-1} 2 k_{2}^{\prime}+1,2^{-j-1}\left(2 k_{2}^{\prime}+2\right)\right),
\end{align*}
$$

as in (22). Then for $l=1,2$, we have that

$$
\begin{align*}
\hat{\mu}_{j}\left(A_{l}\right) & =\mu_{j-1}^{*}\left(A_{l}\right)+\sum_{i=1}^{3} v_{i}\left(A_{l}\right) \\
& =\mathbb{E}\left[\mu_{j-1}^{*} \chi_{A_{l}}\right]+\sum_{i=1}^{3} \mathbb{E}\left[v_{i} \chi_{A_{l}}\right] . \tag{51}
\end{align*}
$$

Now we will calculate each one of the expectations separately. By (17), (18) and (25), we have that

$$
\begin{equation*}
\mu_{j-1}^{*}(x, y)=\sum_{k_{1}} \sum_{k_{2}}\left\langle\mu_{j-1}^{*}, \Phi_{j, k_{1}, k_{2}}\right\rangle \Phi_{j, k_{1}, k_{2}}(x, y) \tag{52}
\end{equation*}
$$

and

$$
\begin{align*}
v_{1}(x, y) & =\sum_{k_{1}} \sum_{k_{2}}\left\langle v_{1}, \Psi_{j, k_{1}, k_{2}}^{(1)}\right\rangle \Psi_{j, k_{1}, k_{2}}^{(1)}(x, y) \\
& =\sum_{k_{1}} \sum_{k_{2}}\left\langle v_{1}, \Psi_{j, k_{1}, k_{2}}^{(1)}\right\rangle\left[\chi_{A_{1}\left(j, k_{1}, k_{2}\right)}(x, y)-\chi_{A_{2}\left(j, k_{1}, k_{2}\right)}(x, y)\right] . \tag{53}
\end{align*}
$$

Then for $l=1$ and by (15), (16) and (22), we have that

$$
\begin{align*}
\mathbb{E}\left[\mu_{j-1}^{*} \chi_{A_{1}}\right] & =\mathbb{E}\left[\chi_{A_{1}} \sum_{k_{1}} \sum_{k_{2}}\left\langle\mu_{j-1}^{*}, \Phi_{j, k_{1}, k_{2}}\right\rangle \Phi_{j, k_{1}, k_{2}}\right] \\
& \left.=\left\langle\mu_{j-1}^{*}, \Phi_{j, k_{1}^{\prime}, k_{2}^{\prime}}\right\rangle \mathbb{E}\left[\chi_{A_{1}} \Phi_{j, k_{1}^{\prime}, k_{2}}\right]\right]  \tag{54}\\
& =d \mathbb{E}\left[\chi_{A_{1}}\right] \\
& =d \frac{1}{2^{2 j+1}}
\end{align*}
$$

where $d=\left\langle\mu_{j-1}^{*}, \Phi_{j, k_{1}^{\prime}, k_{2}^{\prime}}\right\rangle$ and

$$
\begin{align*}
\mathbb{E}\left[v_{1} \chi_{A_{1}}\right] & =\mathbb{E}\left[\chi_{A_{1}} \sum_{k_{1}} \sum_{k_{2}}\left\langle v_{1}, \Psi_{j, k_{1}, k_{2}}^{(1)}\right\rangle\left(\chi_{A_{1}\left(j, k_{1}, k_{2}\right)}-\chi_{A_{2}\left(j, k_{1}, k_{2}\right)}\right)\right] \\
& =\left\langle v_{1}, \Psi_{j, k_{1}^{\prime}, k_{2}^{\prime}}^{(1)}\right\rangle \mathbb{E} \chi_{A_{1}}  \tag{55}\\
& =c \mathbb{E}\left[\chi_{A_{1}}\right] \\
& =c \frac{1}{2^{2 j+1}}
\end{align*}
$$

where $c=\left\langle v_{1}, \Psi_{j, k_{1}^{\prime}, k_{2}^{\prime}}^{(1)}\right\rangle$. Similarly, but using (23) and (24), we can prove that

$$
\begin{equation*}
\mathbb{E}\left[v_{2} \chi_{A_{1}}\right]=0 \tag{56}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbb{E}\left[v_{3} \chi_{A_{1}}\right]=0 \tag{57}
\end{equation*}
$$

Then

$$
\begin{align*}
\hat{\mu}_{j}\left(A_{1}\right) & =d \frac{1}{2^{2 j+1}}+c \frac{1}{2^{2 j+1}}  \tag{58}\\
& =\mu_{j-1}^{*}\left(A_{1}\right)+v\left(A_{1}\right)
\end{align*}
$$

Now, we will make an analogous argument for the case $l=2$. Hence,

$$
\begin{align*}
\mathbb{E}\left[\mu_{j-1}^{*} \chi_{A_{2}}\right] & =\mathbb{E}\left[\chi_{A_{2}} \sum_{k_{1}} \sum_{k_{2}}\left\langle\mu_{j-1}^{*}, \Phi_{j, k_{1}, k_{2}}\right\rangle \Phi_{j, k_{1}, k_{2}}\right] \\
& \left.=\left\langle\mu_{j-1}^{*}, \Phi_{j, k_{1}, k_{2}}\right\rangle \mathbb{E}\left[\chi_{A_{2}} \Phi_{j, k_{1}^{\prime}, k_{2}^{\prime}}\right]\right]  \tag{59}\\
& =d \mathbb{E}\left[\chi_{A_{2}}\right] \\
& =d \frac{1}{2^{2 j+1}}
\end{align*}
$$

and

$$
\begin{align*}
\mathbb{E}\left[v_{1} \chi_{A_{2}}\right] & =\mathbb{E}\left[\chi_{A_{2}} \sum_{k_{1}} \sum_{k_{2}}\left\langle v_{1}, \Psi_{j, k_{1}, k_{2}}^{(1)}\right\rangle\left(\chi_{A_{1}\left(j, k_{1}, k_{2}\right)}-\chi_{A_{2}\left(j, k_{1}, k_{2}\right)}\right)\right] \\
& =-\left\langle v_{1}, \Psi_{j, k_{1}^{\prime}, k_{2}^{\prime}}^{(1)}\right\rangle \mathbb{E} \chi_{A_{2}}  \tag{60}\\
& =-c \mathbb{E}\left[\chi_{A_{2}}\right] \\
& =-c \frac{1}{2^{2 j+1}}
\end{align*}
$$

Thus,

$$
\begin{align*}
\hat{\mu}_{j}\left(A_{2}\right) & =d \frac{1}{2^{2 j+1}}-c \frac{1}{2^{2 j+1}}  \tag{61}\\
& =\mu_{j-1}^{*}\left(A_{2}\right)+v\left(A_{2}\right)
\end{align*}
$$

and by (54), (55), (59) and (60), we have that

$$
\begin{equation*}
\mu_{j-1}^{*}\left(A_{1}\right)=\mu_{j-1}^{*}\left(A_{2}\right) \tag{62}
\end{equation*}
$$

and

$$
\begin{equation*}
v_{1}\left(A_{1}\right)=-v_{1}\left(A_{2}\right) \tag{63}
\end{equation*}
$$

Therefore, it follows from (58), (61), (62), (63) and the fact that $\hat{\mu}_{j}$ is a positive measure, that

$$
\begin{align*}
\left|v_{1}\left(A_{1}\right)\right| & =\left|v_{1}\left(A_{2}\right)\right| \\
& \leq\left|\mu_{j-1}^{*}\left(A_{1}\right)\right|  \tag{64}\\
& =\left|\mu_{j-1}^{*}\left(A_{2}\right)\right|
\end{align*}
$$

From which we conclude that $\mu_{j-1}^{*}(A)=0$ implies $\nu_{1}(A)=0$. Similarly, it can be shown that $\mu_{j-1}^{*}(A)=0$ implies $v_{i}(A)=0$ for $i=2,3$; for this, analogous proofs are carried out, with the difference being that for $i=2$, the sets to be considered are $B_{1}$ and $B_{2}$ as in (23), whereas $C_{1} \cup C_{4}$ and $C_{2} \cup C_{3}$ as in (24) are the respective sets when $i=3$.

We have proved that if it is intended to go back to the preimage of the projection of the approximation operator $\mathbf{P}$ from a level $j-1$ to a level $j$, the support of the level $j-1$ delimits that of the $j$ level. Now, we will prove that for every measure $\hat{\mu}_{j}$ that satisfies (45) and (46), it necessarily holds that $v_{1}=v_{2}=0$.

Lemma 2. Let $j \in \mathbb{Z}$ and $\mu_{j-1}^{*}$ be an optimal simple solution for (40) at level $j-1$. Then for each feasible solution $\hat{\mu}_{j} \in \mathbf{P}_{j} L^{2}\left(\mathbb{R}^{2}\right)$ to (40) at level $j$ such that $\mathbf{P}_{j-1} \hat{\mu}_{j}=\mu_{j-1}^{*}$ and $\hat{\mu}_{j}=$ $\mu_{j-1}^{*}+v_{1}+v_{2}+v_{3}$. It necessarily holds that

$$
\begin{equation*}
v_{1}=v_{2}=0 \tag{65}
\end{equation*}
$$

Proof. In order to perform this proof, we use the restrictions of the MK problem (40), which in turn, are related to the marginal measures. Therefore, we will only complete the proof for one of the projections, since the other is analogous. From the linearity of the Radon-Nikodym derivative, it follows that

$$
\begin{align*}
\eta_{j}^{1} & =\Pi \hat{\mu}_{j} \\
& =\Pi\left(\mu_{j-1}^{*}+v_{1}+v_{2}+v_{3}\right)  \tag{66}\\
& =\Pi \mu_{j-1}^{*}+\Pi v_{1}+\Pi v_{2}+\Pi v_{3} .
\end{align*}
$$

Let $k^{\prime} \in \mathbb{Z}$ and $I_{j+1,2 k^{\prime}}=\left[2^{-j-1} 2 k^{\prime}, 2^{-j-1}\left(2 k^{\prime}+1\right)\right)$. Thus, by (39) and (66), we have that

$$
\begin{align*}
\Pi_{1} \hat{\mu}_{j}\left(I_{j+1,2 k^{\prime}}\right) & =\eta_{j}^{1}\left(I_{j+1,2 k^{\prime}}\right) \\
& =\Pi_{1} \mu_{j-1}^{*}\left(I_{j+1,2 k^{\prime}}\right)  \tag{67}\\
& =\alpha .
\end{align*}
$$

That is, we are evaluating feasible solutions on rectangles whose height is half the size of the squares with which they are discretized at the $j$ level of the Haar MRA. Now, we will develop in detail (66) evaluated on $I_{j, k^{\prime}}$. Since $\mu_{j-1}^{*}$ is a simple solution, we call $l^{\prime} \in \mathbb{Z}$ the only number such that $\mu_{j-1}^{+}\left(S\left(j, k^{\prime}, l^{\prime}\right)\right)>0$, where $S\left(j, k_{1}, k_{2}\right)$ is defined as in (20). With the aim of simplifying the notation, we define $\mathcal{I}=I_{j+1,2 k^{\prime}} \times \mathbb{R}$. By (16) and (25), we have that

$$
\begin{align*}
\Pi_{1} v_{1}\left(I_{j+1,2 k^{\prime}}\right) & =v_{i}(\mathcal{I}) \\
& =\mathbb{E}\left[\chi_{\mathcal{I}} v_{1}\right] \\
& =\mathbb{E}\left[\chi_{\mathcal{I}} \sum_{k_{1}} \sum_{k_{2}}\left\langle v_{1}, \Psi_{j, k_{1}, k_{2}}^{(1)}\right\rangle\left(\chi_{A_{1}\left(j, k_{1}, k_{2}\right)}-\chi_{A_{2}\left(j, k_{1}, k_{2}\right)}\right)\right]  \tag{68}\\
& =\mathbb{E}\left[\chi_{\mathcal{I}} \sum_{k_{2}}\left\langle v_{1}, \Psi_{j, k^{\prime}, k_{2}}^{(1)}\right\rangle\left(\chi_{A_{1}\left(j, k^{\prime}, k_{2}\right)}-\chi_{A_{2}\left(j, k^{\prime}, k_{2}\right)}\right)\right] \\
& =\mathbb{E}\left[\chi_{\mathcal{I}}\left\langle v_{1}, \Psi_{j, k^{\prime}, l^{\prime}}^{(1)}\right\rangle\left(\chi_{A_{1}\left(j, k^{\prime}, l^{\prime}\right)}-\chi_{A_{2}\left(j, k^{\prime}, l^{\prime}\right)}\right)\right]
\end{align*}
$$

By the way $\mathcal{I}$ was defined, necessarily in the last equality it must be fulfilled that one of the terms in the expectation is equal to 0 . Hence, it is fulfilled that

$$
\begin{align*}
\Pi_{1} v_{1}\left(I_{j+1,2 k^{\prime}}\right) & =\left\langle v_{1}, \Psi_{j, k^{\prime}, l^{\prime}}^{(1)}\right\rangle \mathbb{E}\left[\chi_{\mathcal{I}}\right] \\
& =\frac{1}{2^{2 j+1}}\left\langle v_{1}, \Psi_{j, k^{\prime}, l^{\prime}}^{(1)}\right\rangle . \tag{69}
\end{align*}
$$

By a similar argument, it can be proved that

$$
\begin{equation*}
\Pi v_{2}\left(I_{j+1,2 k^{\prime}}\right)=0 \tag{70}
\end{equation*}
$$

and

$$
\begin{equation*}
\Pi v_{2}\left(I_{j+1,2 k^{\prime}}\right)=0 . \tag{71}
\end{equation*}
$$

Then from (66) to (71), it follows that

$$
\begin{equation*}
\alpha=\alpha+\frac{1}{2^{2 j+1}}\left\langle v_{1}, \Psi_{j, k^{\prime}, l^{\prime}}^{(1)}\right\rangle+0+0 . \tag{72}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
\left\langle v_{1}, \Psi_{j, k^{\prime}, l^{\prime}}^{(1)}\right\rangle=0 . \tag{73}
\end{equation*}
$$

Therefore, $v_{1}=0$. In a similar way, we can prove that $v_{2}=0$.

Suppose we have a simple optimal solution $\mu_{j-1}^{*}$ for the MK problem discretized through the Haar MRA at level $j-1 \in \mathbb{Z}$ and that we are interested in refining that solution to the next level $j$. By Lemmas 1 and 2 , any $\hat{\mu}_{j}$ that satisfies (45) has its support contained in the support of $\mu_{j-1}^{*}$ and has components $v_{1}=v_{2}=0$. Then the problem of constructing a feasible solution $\hat{\mu}_{j}$ such that it refines $\mu_{j-1}^{*}$ is reduced to construct

$$
\begin{align*}
\hat{\mu}_{j}(x, y) & =\mu_{j-1}^{*}(x, y)+v_{3}(x, y) \\
& =\mu_{j-1}^{*}(x, y)+\sum_{k_{1}} \sum_{k_{2}}\left\langle\hat{\mu}_{j}, \Psi_{j, k_{1}, k_{2}}^{(i)}\right\rangle \Psi_{j, k_{1}, k_{2}}^{(i)}(x, y), \tag{74}
\end{align*}
$$

which is equivalent to chosing for each $k_{1}, k_{2} \in \mathbb{Z}$ a value $v_{k_{1}, k_{2}}^{3}$ such that

$$
\begin{equation*}
\hat{\mu}_{k_{1}, k_{2}}^{j}=\mu_{k_{1}, k_{2}}^{*, j-1}+v_{k_{1}, k_{2}}^{3} . \tag{75}
\end{equation*}
$$

By Lemma 1 for each $k_{1}, k_{2} \in \mathbb{Z}$, it is fulfilled that

$$
\begin{equation*}
\left|v_{k_{1}, k_{2}}^{3}\right| \leq \mu_{k_{1}, k_{2}}^{*, j-1} \tag{76}
\end{equation*}
$$

Therefore, the choice of $v_{k_{1}, k_{2}}^{3}$ is restricted to a compact collection, and since $\hat{\mu}_{j}$ is a solution of the linear program (40), then

$$
\begin{equation*}
\left|v_{k_{1}, k_{2}}^{3}\right|=\mu_{k_{1}, k_{2}}^{*, j-1} \tag{77}
\end{equation*}
$$

Thus, the sign of $v_{k_{1}, k_{2}}^{3}$ must be such that it minimizes $v_{k_{1}, k_{2}}^{3} c_{k_{1}, k_{2}}^{j}$. That is,

$$
\begin{equation*}
v_{k_{1}, k_{2}}^{3} c_{k_{1}, k_{2}}^{j} \leq 0 \tag{78}
\end{equation*}
$$

## 4. Methodological Proposal

In this section, we show through examples a process that builds solutions to the MK problem with a reduced number of operations. First, we consider the MK problem with cost function $c:[0,1] \times[0,1] \rightarrow \mathbb{R}$ defined by

$$
\begin{equation*}
c(x, y)=4 x^{2} y-x y^{2} \tag{79}
\end{equation*}
$$

and homogeneous restrictions $v_{1}, v_{2}$ over $[0,1]$. So that the algorithm can be graphically appreciated, we take a small level of discretization, namely $j=6$. Thus, in the Haar MRA over $\mathbb{R}^{2}$ at level $j=6$, the cost function has the form shown in Figure 5, which can be stored in a vector of size $2^{2 j}=2^{2(6)}=2^{12}$.

Now, we apply the filtering process to the cost function at level $j=6$, which results in four functions

$$
\begin{equation*}
c_{6}(x, y)=c_{5}(x, y)+\Psi^{(1)}(x, y)+\Psi^{(2)}(x, y)+\Psi^{(3)}(x, y) ; \tag{80}
\end{equation*}
$$

see Figures 6-9.


Figure 5. Step 1. Discretization of the cost function to the level $j$, which is denoted by $c_{j}$. In particular, the cost function is $c(x, y)=4 x^{2} y-x y^{2}$ for lever $j=6$.

Step 2. Filtering the original discrete function using the high-pass filter, which yield three discrete functions denoted by $\Psi_{1}, \Psi_{2}$ and $\Psi_{3}$, that functions correspond to Figures 7, 8 and 9 , respectively, each describing local changes in the original discrete function. It is then low-pass filtered to produce the approximate discrete function $c_{5}$, which is given by Figure 6.


Figure 6. $c_{5}(x, y)$.


Figure 7. $\Psi^{(1)}(x, y)$.


Figure 8. $\Psi^{(2)}(x, y)$.


Figure 9. $\Psi^{(3)}(x, y)$.
We then solve the MK problem for the level $j-1=5$. That is, we find a measure $\mu_{j-1}^{*}=\mu_{5}^{*}$ that is an optimal solution for the MK problem with cost function $c_{5}$. Such data can be stored in a vector of size $2^{2 j-2}=2^{10}$; see Figure 10. For each entry $k$, the formal application that plots this vector in a square is defined by $k \rightarrow\left(k_{1}+1, k_{2}+1\right)$, where $0 \leq k_{1}, k_{2}<2^{n-1}$ and $k=k_{1} \cdot 2^{j-1}+k_{2}$.


Figure 10. Step 3. We obtain a solution $\mu_{5}^{*}$ for $M K^{5}$ associated with the cost function $c_{5}$ given in Figure 6.

Since the measure $\mu_{j-1}^{*}=\mu_{5}^{*}$ is an optimal simple solution for the MK problem, then we can represent its support in a simple way, as we show below:

$$
\begin{equation*}
\operatorname{support}\left(\mu_{j-1}^{*}\right)=\bigcup_{k_{1}, k_{2}} S_{k_{1}, k_{2}} \tag{81}
\end{equation*}
$$

where $S_{k_{1}, k_{2}}=S\left(j, k_{1}, k_{2}\right)$ is the square in (20). Next, we split each block $S_{k_{1}, k_{2}} \subset \operatorname{support}\left(\mu_{j-1}^{*}\right)$ into four parts as in (24); see Figure 11.


Figure 11. Division of the components of $\operatorname{support}\left(\mu_{j-1}^{*}\right)$ into four parts.
From the technical point of view, in the discretization at level $j-1=5$, we have a grid of $2^{j-1} \times 2^{j-1}=32 \times 32$ squares that we call $S_{k_{1}, k_{2}}$ and identify with the points $\left(k_{1}, k_{2}\right)$. Thus, we refine to a grid of $64 \times 64$, splitting each square into four, which in the new grid are determined by points

$$
\begin{equation*}
\left(2 k_{1}-1,2 k_{2}\right),\left(2 k_{1}-1,2 k_{2}-1\right),\left(2 k_{1}, 2 k_{2}-1\right) \text { and }\left(2 k_{1}, 2 k_{2}\right) ; \tag{82}
\end{equation*}
$$

see Figure 12.


Figure 12. Refinement of grid from level $j-1$ to $j$ of discretization.
As we prove in Lemma 1, any feasible solution $\hat{\mu}_{j}=\hat{\mu}_{6}$ that refines $\mu_{5}^{*}$ has its support contained in $\operatorname{support}\left(\mu_{5}^{*}\right)$. Therefore, we must only deal with the region delimited by the support of $\mu_{5}^{*}$. By Lemma 2 , in order to construct the solution $\hat{\mu}_{6}$, we only need to determine the values $n_{k_{1}, k_{2}}^{3}$ corresponding to the coefficients of the wavelet part $v_{3}$; however, by (77) and (78), those values are well determined and satisfy that when added with the scaling part $\mu_{k_{1}, k_{2}}^{*}$, the result is a scalar multiple of a characteristic function. For example if the
square $S_{k_{1}, k_{2}}$ has scaling part coefficient $\mu_{k_{1}, k_{2}}^{*}=c$, then we choose $v_{k_{1}, k_{2}}^{3}=-c$. Hence, by (21) and (25), we have that

$$
\begin{align*}
\hat{\mu}_{k_{1}, k_{2}}^{6}(x, y) & =\mu_{k_{1}, k_{2}}^{*}(x, y)+v_{k_{1}, k_{2}}^{3}(x, y) \\
& =c \chi_{S}(x, y)-c\left[\chi_{C_{1} \cup C_{4}}(x, y)-\chi_{C_{2} \cup C_{3}}(x, y)\right]  \tag{83}\\
& =2 c \chi_{C_{2} \cup C_{3}}(x, y) .
\end{align*}
$$

Thus, from an operational point of view, we only need to chose between two options of supporting each division of square $S_{k_{1}, k_{2}}$, as we illustrate in Figures 13 and 14.


Figure 13. Supports for refinement of the element $S_{k_{1}, k_{2}}$ corresponding to the level $j$ : Option I.


Figure 14. Supports for refinement of the element $S_{k_{1}, k_{2}}$ corresponding to the level $j$ : Option II.
This coincides with our geometric intuition. Hence, the resulting feasible solution $\hat{\mu}_{6}$ has its support contained in $\operatorname{support}\left(\mu_{5}^{*}\right)$, and its weight within each square $S_{k_{1}, k_{2}}$ is presented in a diagonal within that block; see Figure 15. Finally, we can improve $\hat{\mu}_{6}$ by observing the way the filtering process acts. To do this, we apply the proximity criteria (41) and split the support of $\hat{\mu}_{6}$ into points of Type I and II. In Figure 16, we identify the points of Type I and II of solution $\mu_{5}^{*}$, whereas in Figure 15, we do the same but for $\hat{\mu}_{6}$.

Intuitively, the division of the support into points of Type I and II allows us to classify the points so that they have an identity function form and, consequently, that come from the discretization of a continuous function-points of Type I—and in points that come from the discretization of a discontinuous function-points of Type II. Thus, the points of Type I are located in such a way that they generate a desired solution, and therefore it is not convenient to move them, whereas Type II points are free to be changed as this does not lead to the destruction of a continuous structure in the solution. As we mentioned in the previous section, each permutation of rows or columns of one weighted element $S_{k_{1}, k_{2}}$ with another $S_{k_{1}^{\prime}, k_{2}^{\prime}}$ constructs a feasible solution; see (44). Thus, as a heuristic technique to improve the solution, we check the values $\left\langle c_{6}, \mu_{6}^{\sigma}\right\rangle$ associated with each solution $\mu_{6}^{\sigma}$ obtained by permuting rows or columns of points of Type II of the solution $\hat{\mu}_{6}$. We call $\mu_{6}^{\sigma^{*}}$ the solution for which its permutation gives it the best performance. See Figure 17.


Figure 15. Step 5. Using the solution $\mu_{5}^{*}$ which is given by Figure 10, the information of the high-pass components (Figures 7, 8 and 9) and Lemma 1, we obtain a feasible solution for Level 6, which is denoted by $\hat{\mu}_{6}$.


Figure 16. Step 4. We classify the points of the support of the solution $\mu_{5}^{*}$ by proximity criteria as points of Type I $\square$ or Type II $\square$ (the measure $\mu_{5}^{*}$ corresponds to Figure 10).

Finally, we present Table 1 that compare the solutions of the MK problem, in which $\mathrm{MK}^{5}$ is the value associated with the optimal solution $\mu_{5}^{*}$ at level of discretization $j-1=5$, $\mathrm{MK}_{6}$ is the value associated with an optimal solution $\mu_{6}^{*}$ at level of discretization $j=6$, and $\mathrm{MK}_{6}^{\sigma^{*}}$ is the value associated with the solution $\mu_{6}^{\sigma^{*}}$ obtained by the heuristic method described in the previous paragraph.


Figure 17. Step 6. Classification of the points of $\mu_{5}^{*}$ induces classification of the points in $\hat{\mu}_{6}$ by contention in the support. Over the points of Type I of the solution $\hat{\mu}_{6}$, we do not move those points. For the points of Type II, we apply a permutation to the solution over the two points that improve the solution and repeat the process with the rest of the points.

Table 1. Comparison of the values corresponding to $\mathrm{MK}_{5}, \mathrm{MK}_{6}^{\sigma^{*}}$ and $\mathrm{MK}_{6}$.

| $\mathbf{M K}_{\mathbf{5}}$ | MK $_{6}^{\sigma^{*}}$ | MK $_{6}$ |
| :---: | :---: | :---: |
| 0.24785 | 0.247738 | 0.247726 |

## 5. Other Examples of This Methodology

We conclude this work with a series of examples in which we apply the proposed methodology. Each of them is divided into the six-step algorithm introduced in the previous section and corresponds to a classical example existing in the literature.

### 5.1. Example with Cost Function $c(x, y)=x^{2} y-x y^{2}$

Let MK be the MK problem with cost function $c(x, y)=x^{2} y-x y^{2}$ and uniform restrictions $\eta_{1}$ and $\eta_{2}$ over the unitary interval $[0,1]$. In order to be more didactic, we consider the level of discretization $j=6$. Next, we present the reduced algorithm.

Step 1. Discretize the cost function at level $j$-that is, over a grid of $2^{j} \times 2^{j}$. Figure 18.


Figure 18. Discretization of the cost function $c(x, y)=x^{2} y-x y^{2}$ at level $j=6$.

Step 2. Apply the filtering process to the cost function at level $j-1$, obtaining the filtered cost function $c_{5}$, which is plotted on a grid of $2^{j-1} \times 2^{j-1}$. Figure 19.


Figure 19. Filtering of the cost function $c$ at level $j-1=5$.
Step 3. Find an optimal simple solution $\mu_{j-1}^{*}$ for the discretized $M K_{j}$ problem. That is, solve for the cost function $c_{5}$ and obtain an optimal simple solution $\mu_{5}^{*}$. Figure 20.


Figure 20. MK solution for the filtered function $c$ at level $j-1=5$.
Step 4. Apply the proximity criteria to the support of $\mu_{j-1}^{*}$. Figure 21.


Figure 21. Classification of points of $\operatorname{support}\left(\mu_{5}^{*}\right)$ into Type I $\square$ and Type II $\square$.

Step 5. Refine the optimal simple solution $\mu_{j-1}^{*}$ to a feasible solution $\hat{\mu}_{j}$, dividing each weighted square $S_{k_{1}, k_{2}}$ at level $j-1$ into four squares at level $j$ (see (82)) and place mass according to the criteria (83). Figure 22.


Figure 22. Solution $\hat{\mu}_{6}$ from refinement of $\mu_{5}^{*}$.
Step 6. Permute the rows and columns of the points of Type II in support $\left(\hat{\mu}_{j}\right)$ using (44) to construct feasible solutions $\mu_{j}^{\sigma}$ and chose the one that has better performance. Figure 23.


Figure 23. Final result $\mu_{6}^{\sigma^{*}}$.
The following Table 2 contains the comparison of the proposed methodology with the two immediate levels of resolution.

Table 2. Comparison of the values corresponding to $\mathrm{MK}_{5}, \mathrm{MK}_{6}^{\sigma^{*}}$ and $\mathrm{MK}_{6}$.

| $\mathbf{M K}_{\mathbf{5}}$ | $\mathbf{M K}_{\mathbf{6}}{ }^{\boldsymbol{*}}$ | $\mathbf{M K}_{\mathbf{6}}$ |
| :---: | :---: | :---: |
| -0.0350952 | -0.035141 | -0.035141 |

### 5.2. Example with Cost Function $c(x, y)=x y$

Let MK be the MK problem with cost function $c(x, y)=x y$ and uniform restrictions $\eta_{1}$ and $\eta_{2}$ over the unitary interval $[0,1]$. In order to be more didactic, we consider the level of discretization $j=6$. Next, we present the reduced algorithm.
Step 1. Discretize the cost function at level $j$-that is, over a grid of $2^{j} \times 2^{j}$. Figure 24.


Figure 24. Discretization of the cost function $c(x, y)=x y$ at level $j=6$.
Step 2. Apply the filtering process to the cost function at level $j-1$, obtaining the filtered cost function $c_{5}$, which is plotted on a grid of $2^{j-1} \times 2^{j-1}$. Figure 25 .


Figure 25. Apply the filter to the cost function $c$ at level $j-1=5$.
Step 3. Find an optimal simple solution $\mu_{j-1}^{*}$ for the discretized $M K_{j}$ problem. That is, solve for the cost function $c_{5}$ and obtain an optimal simple solution $\mu_{5}^{*}$. Figure 26.


Figure 26. Solution of the MK for the function $c$ at level $j-1=5$.
Step 4. Apply the proximity criteria to the support of $\mu_{j-1}^{*}$. Figure 27.


Figure 27. In this example there are only Type I points
Step 5. Refine the optimal simple solution $\mu_{j-1}^{*}$ to a feasible solution $\hat{\mu}_{j}$, dividing each weighted square $S_{k_{1}, k_{2}}$ at level $j-1$ into four squares at level $j$ (see (82)) and placing mass according to the criteria (83). Figure 28.


Figure 28. Solution $\hat{\mu}_{6}$ from refinement of $\mu_{5}^{*}$.
Step 6. Permute the rows and columns of the points of Type II in support $\left(\hat{\mu}_{j}\right)$ using (44) to construct feasible solutions $\mu_{j}^{\sigma}$ and chose the one that has better performance. Figure 29.


Figure 29. Final result $\mu_{6}^{\sigma^{*}}$.

The following Table 3 contains the comparison of the proposed methodology with the two immediate levels of resolution.

Table 3. Comparison of the values corresponding to $\mathrm{MK}_{5}, \mathrm{MK}_{6}^{\sigma^{*}}$ and $\mathrm{MK}_{6}$.

| MK $_{\mathbf{5}}$ | MK $_{6}^{\sigma^{*}}$ | MK $_{\mathbf{6}}$ |
| :---: | :---: | :---: |
| 0.166748 | 0.166687 | 0.166687 |

### 5.3. Example with Cost Function $c(x, y)=(2 y-x-1)^{2}(2 y-x)^{2}$

We take the MK problem with cost function $c(x, y)=(2 y-x-1)^{2}(2 y-x)^{2}$ and homogeneous restrictions over the unitary interval. Again, we consider the level of discretization $j=6$.
Step 1. Discretize the cost function at level $j$. Figure 30.


Figure 30. Discretization of cost function $c(x, y)=(2 y-x-1)^{2}(2 y-x)^{2}$ for level $j=6$.
Step 2. Apply the filtering process at level $j-1$ to the cost function. Figure 31.


Figure 31. Filtered cost function for $j=5$.
Step 3. Find an optimal simple solution $\mu_{j-1}^{*}$ for the $\mathrm{MK}_{j-1}$ problem. Figure 32.


Figure 32. Solution of the $M K_{5}$ problem.
Step 4. Apply the proximity criteria to the support of $\mu_{j-1}^{*}$. Figure 33.


Figure 33. In this example, there are only Type II points
Step 5. Refine the solution $\mu_{j-1}^{*}$ to a feasible solution $\hat{\mu}_{6}$ applying the criteria (82) and (83). Figure 34.


Figure 34. Feasible solution $\hat{\mu}_{6}$ from refinement of $\mu_{5}^{*}$.
Step 6. Permute points of Type II of $\hat{\mu}_{j}$ using (44) to construct feasible solutions $\mu_{j}^{\sigma}$ and chose which has better performance. Figure 35.


Figure 35. Final result $\mu_{6}^{\sigma^{*}}$.
The Table 4 summarizes the results obtained.
Table 4. Comparison of values $\mathrm{MK}_{5}, \mathrm{MK}_{6}^{\sigma^{*}}$ and $\mathrm{MK}_{6}$.

| $\mathbf{M K}_{\mathbf{5}}$ | $\mathbf{M K}_{\mathbf{6}}{ }^{\boldsymbol{*}}$ | $\mathbf{M K}_{\mathbf{6}}$ |
| :---: | :---: | :---: |
| 0.000236571 | 0.0000600852 | 0.0000597889 |

## 6. Conclusions and Future Work

Note that with the methodology of [3,4,9], the authors obtain a solution of $M K^{j}$. For this, they need to resolve a transport problem with $4^{j}$ variables. We call this methodology an exhaustive method. For our methodology, in Step 3, we need to resolve a transport problem with $4^{j-1}$ variables, and the other steps of the methodology are methods of classification, ordering and filtering; with $2^{j}$ data for classification and ordering and $4^{j}$ data for filtering, it is clear that this method requires fewer operations to resolve transport problems.

In summary, we have the following table comparing the results of solving the examples more often used in the literature with our methodology versus the exhaustive method (using all variables).

| Cost Function | $M K^{6}$ | $M K_{6}^{\sigma}$ | Error |
| :---: | :---: | :---: | :---: |
| $x y$ | 0.166687 | 0.166687 |  |
| $x^{2} y-x y^{2}$ | -0.035141 | -0.035141 | $1.52588 \times 10^{-5}$ |
| $4 x^{2} y-x y^{2}$ | 0.247726 | 0.247738 | $5.31762 \times 10^{-5}$ |
| $(2 y-x-1)^{2}(2 y-x)^{2}$ | 0.0000600947 | 0.0000600947 |  |
| Cost Function | $M K^{7}$ | $M K_{7}^{\sigma}$ | Error |
| $x y$ | 0.166672 | 0.166672 |  |
| $x^{2} y-x y^{2}$ | -0.0351524 | -0.0351524 | $3.8147 \times 10^{-6}$ |
| $4 x^{2} y-x y^{2}$ | 0.247695 | 0.247698 | $1.24066 \times 10^{-5}$ |
| $(2 y-x-1)^{2}(2 y-x)^{2}$ | 0.0000600852 | 0.0000600852 |  |
| Cost Function | $M K^{8}$ | $M K_{8}^{\sigma}$ | Error |
| $x y$ | 0.166668 | 0.166668 |  |
| $x^{2} y-x y^{2}$ | -0.0351553 | -0.0351553 | $9.53674 \times 10^{-7}$ |
| $4 x^{2} y-x y^{2}$ | 0.247688 | 0.247688 | $3.21631 \times 10^{-6}$ |
| $(2 y-x-1)^{2}(2 y-x)^{2}$ |  |  |  |

Note that our method always improves the solution of the level $j-1$ and for some examples give an exact solution; we use Mathematica ${ }^{\circledR}$ and basic computer equipment
for programming this methodology, and maybe we can improve the results with software focused on numerical calculus and better computer equipment. It is also important to mention that the methodology presented in this work has some weaknesses. In our computational experiments, we noticed that if we did not start with a sufficient amount of information, then the methodology tended to give very distant results. In other words, if the initial level of discretization was not fine enough, then because the algorithm lowers the resolution level when executed, such loss of information generates poor performance. However, when starting with an adequate level of discretization, experimentally it can be observed that the distribution of the solutions for the discretized problems, as well as the respective optimal values, have stable behavior with a clear trend. The question that arises naturally is: "In practice, what are the parameters that determine good or bad behavior of the algorithm?" Clearly, if the cost function is fixed and we rule out the possible technical problems associated with programming and computing power, the only remaining parameter is the initial refinement level at which the algorithm is going to work-that is, the level $j$. However, if we reflect more deeply on the reasons why there is a practical threshold beyond which at a certain level of discretization the algorithm has stable behavior, we only have as a possible causes the level of information of the cost function that captures the MR analysis. In other words, if the oscillation of the cost function at a certain level of resolution is well determined by MR analysis, then the algorithm will have good performance.

The approach presented in this paper is far from exhaustive and, on the contrary, opens the possibility for a number of new proposals for approximating solutions to the MK problem. The above is due to the fact that in the work [9], it was proven that discretization of the MK problem can be performed from any MR analysis over $\mathbb{R}^{2}$. Therefore, the possibility of implementing other types of discretions remains open. In principle, as we mentioned in the previous paragraph, the most natural thing is to expect better performance if the nature of the cost function and the types of symmetric geometric structures that it induces in space are studied in order to use an MR analysis that fits this information and therefore has more efficient performance.

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# A License Plate Recognition System with Robustness against Adverse Environmental Conditions Using Hopfield's Neural Network 

Saman Rajebi ${ }^{1}$, Siamak Pedrammehr ${ }^{2}$ (D) and Reza Mohajerpoor ${ }^{3, *}$<br>1 Department of Electrical Engineering, Seraj University, Tabriz 5137894797, Iran<br>${ }^{2}$ Faculty of Design, Tabriz Islamic Art University, Tabriz 5164736931, Iran<br>3 School of Civil Engineering, The University of Sydney, Sydney 2006, Australia<br>* Correspondence: reza.mohajerpoor@sydney.edu.au

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#### Abstract

License plates typically have unique color, size, and shape characteristics in each country. This paper presents a general method for character extraction and pattern matching in license plate recognition systems. The proposed method is based on a combination of morphological operations and edge detection techniques, along with the bounding box method for identifying and revealing license plate characters while removing unwanted artifacts such as dust and fog. The mathematical model of foggy images is presented and the sum of gradients of the image, which represents the visibility of the image, is improved. Previous works on license plate recognition have utilized nonintelligent pattern matching techniques. The proposed technique can be applied in a variety of settings, including traffic monitoring, parking management, and law enforcement, among others. The applied algorithm, unlike SOTA-based methods, does not need a huge set of training data and is implemented only by applying standard templates. The main advantages of the proposed algorithm are the lack of a need for a training set, the high speed of the training process, the ability to respond to different standards, the high response speed, and higher accuracy compared to similar tasks.


Keywords: pattern recognition; image processing; independent component analysis; Hopfield's neural network; license plate; LPR

MSC: 68T07

## 1. Introduction

Nowadays, the use of surveillance-based security systems has become increasingly important in various applications, such as home security and traffic monitoring. Object detection is one of the fundamental building blocks of automated surveillance systems. Among the most used techniques for object recognition in surveillance systems is the recognition of vehicle license plates [1]. Automatic license plate recognition is an image-processing-based method that is used for security applications such as controlling access to restricted areas and tracking vehicles.

In real-world applications, simple License Plate Recognition (LPR) systems have low detection accuracy [2]. On the one hand, the effects of external factors such as sunlight and car headlights, license plates with inappropriate designs, the wide variety of license plates and, on the other hand, the limited quality of the software and hardware related to the camera, have reduced the accuracy of these systems. However, recent advances in software and hardware have made LPR systems much safer and more widespread [3,4]. A countless number of these systems are working around the world and are growing exponentially and can do more tasks automatically in different market segments. Even if the recognition is not $100 \%$, a results-dependent side program can compensate for the errors and provide an almost flawless system. For example, to calculate the car's parking
time, from entering to leaving the parking lot, this side program can ignore some ignorable errors in the two recognitions. This intelligent integration can overcome the shortcomings of LPR and produce reliable and fully automated systems [5-7].

Figure 1 shows a typical configuration of an LPR system. The license plate reader software is a Windows background program on a PC and an interface between a set of cameras. The program receives the images of the cameras, and by processing them it extracts the license plates of the cars in traffic. The program then displays the results, and can also send them to other parts of the system such as a camera or LED display via serial communication. It then sends this information to the local database or external databases (through the network).


Figure 1. Typical configuration of an LPR system.
The first step in recognizing a car license plate is to distinguish the car from other objects in the image. For this, the methods presented in previous works can be used [8-12]. In similar works, the use of a convolutional neural network (CNN) replaces parts of the proposed method in this paper. Despite the ease of use of this new neural network, there are major disadvantages associated with CNNs. The main disadvantage is that CNNs take a much longer time to train. Another important disadvantage is the need for larger training datasets (i.e., hundreds or thousands of images), and for their proper annotation, which is a delicate procedure that must be performed by domain experts. Other disadvantages include problems that might occur when using pretrained models on similar and smaller datasets (i.e., of a few hundreds of images, or smaller), optimization issues due to model complexities, and hardware restrictions [13]. However, the proposed algorithm assumes there are only cars on highways or an absence of objects corresponding to license plates in non-car elements (such as humans, etc.). This contribution can overcome the burdens that were present in previous works.

In the literature, the pixel-by-pixel comparison method has been used to match the segments extracted from the image with the defined standard characters. This method,
in addition to the low classification accuracy, lack of identification, and removal of noise, also requires a lot of execution time. We employed the Hopfield neural network [14] to simultaneously speed up the program's execution and increase the precision, while removing noises on the segments extracted from the image. License plate recognition systems are usually used outdoors. The presence of air pollution, fog, and other factors causes the car license plate images to become blurred. By using frequency domain techniques, it is possible to remove the side- and destructive effects of the environment on the image. In this paper, first, the location of the license plate is identified, and then the disturbing effects of the environment are removed. After extracting the license plate image's segments, the Hopfield neural network classifies these segments to corresponding defined characters.

## 2. Methodology and Simulation Results

Figure 2 shows the five main steps of license plate recognition. In this structure, the steps for image scheduling, camera settings, and the saving and transferring of results have been ignored. In the following, each of these five steps will be explained in full detail.


Figure 2. The proposed five main steps of license plate recognition.

### 2.1. Preprocessing of the Image

In tasks based on image processing, such as car license plate recognition [15-18] or eye tracking, etc., the first step is usually to determine the approximate location of the target object. For this, a large number of different cars with different license plate locations were studied. All the studied cars were photographed at the same distance and angle. In all these photos, the place of installation of the license plates was marked. Figure 3 shows the border of the area where it was possible for a license plate to exist, taking a suitable tolerance.


Figure 3. The border of the area where it is possible for a license plate to exist.
The area outside the border of the license plate's zone (which has non-useful information) will be affected by the Blur filter. Applying this filter reduces the calculations and the possibility of errors in future processing. Figure 4 shows a typical image affected by this filter, where non-useful areas have been blurred.


Figure 4. A typical image affected by Blur filter; non-useful areas have been blurred.

### 2.2. Elimination of Adverse Environmental Effects

Before determining the exact location of the license plate, and then its characters' segments, the adverse effects of the environment, such as the effects of possible fog or smoke in the space between the camera and the license plate, should be corrected as much as possible. Equation (1) shows a blurry image relation [19-24]:

$$
\begin{equation*}
I(x)=J(x) t(x)+A(1-t(x)) \tag{1}
\end{equation*}
$$

where $I$ is the intensity of the light in the image, $J$ is the illumination of the scene, $A$ is the general light of the environment, and $t$ is a parameter that describes the part of the light that was not scattered and reached the camera. The elimination of adverse environmental effects means recovering $J, A$, and $t$ from $I$. The term $J(x) t(x)$ in this equation is called direct attenuation, which describes the brightness of the scene and its decay in the environment. The term $A(1-t(x))$ is called ambient light, which comes from the previously scattered
light and leads to a change in the color of the environment. When the space is homogeneous, the transfer coefficient $t$ is described as follows:

$$
\begin{equation*}
t(x)=e^{-\beta d(x)} \tag{2}
\end{equation*}
$$

where $\beta$ is the dispersion coefficient. Equation (2) clearly shows that the image brightness decreases exponentially with its depth $d$.

Equation (1) shows that in the RGB color space, vectors $A, I(x)$, and $J(x)$ are coplanar while their endpoints are located on a single line. The transfer coefficient $t$ can be expressed by

$$
\begin{equation*}
t(x)=\frac{\|A-I(x)\|}{\|A-J(x)\|}=\frac{A^{c}-I^{c}(x)}{A^{c}-J^{c}(x)} \tag{3}
\end{equation*}
$$

where $c$ represents the index of the color channel. In blurred images, $t$ is less than one. Thus, the resolution of the image, which is the sum of the image gradients, is low. The following illustrates this reduction:

$$
\begin{equation*}
\sum_{t}\|\nabla I(x)\|=t \sum_{t}\|\nabla J(x)\|<\sum_{t}\|\nabla J(x)\| \tag{4}
\end{equation*}
$$

The transmission coefficient, $t$, is estimated by maximizing the image resolution, while the intensity $J(x)$ is less than the intensity $A$. The dark channel for a haze-free outer space image is defined as the following: in a non-sky image, at least one color channel has very low brightness in some pixels. In other words, the image brightness in these pixels is minimum. Equation (5) shows the dark channel definition:

$$
\begin{equation*}
J^{\text {dark }}(x)=\min _{c \in\{r \cdot g . b\}}\left(\min _{y \in \Omega(x)}\left(J^{c}(y)\right)\right) \tag{5}
\end{equation*}
$$

where $J^{c}$ is a color channel of $J$ and $\Omega(x)$ is a piece of the image centered at $x$. If the image does not include the sky and does not have fog, the intensity of $J^{\text {dark }}$ will be almost zero. Assuming the value of $A$ for the ambient light and the constant transmission coefficient $t(x)$ in a piece of the image, minimizing the intensity (1) gives:

$$
\begin{equation*}
\min _{y \in \Omega(x)}\left(I^{c}(x)\right)=t(x) \min _{y \in \Omega(x)}\left(J^{c}(y)\right)+A^{c}(1-t(x)) \tag{6}
\end{equation*}
$$

Dividing the sides of Equation (6) by $A$ and minimizing again, this time among the color channels gives:

$$
\begin{equation*}
\min _{c}\left(\min _{y \in \Omega(x)}\left(\frac{I^{c}(x)}{A^{c}}\right)\right)=t(x) \min _{c}\left(\min _{y \in \Omega(x)}\left(\frac{J^{c}(y)}{A^{c}}\right)\right)+(1-t(x)) \tag{7}
\end{equation*}
$$

By approximating Equation (7) to zero, the transfer coefficient can be defined as Equation (8). The coefficient $w$ is defined to adjust the blurring in the image.

$$
\begin{equation*}
t(x)=1-w \min _{c}\left(\min _{y \in \Omega(x)}\left(\frac{I^{c}(x)}{A^{c}}\right)\right) \tag{8}
\end{equation*}
$$

By limiting the transmission coefficient on the limit of $t_{0}$, the brightness of the image is expressed as:

$$
\begin{equation*}
J(x)=\frac{I(x)-A}{\max \left(t(x), t_{0}\right)}+A \tag{9}
\end{equation*}
$$

According to the above-stated contents and the presented equations, using the algorithm shown in Figure 5 it is possible to reduce the image blurring to an acceptable level.


Figure 5. Implemented algorithm for removing or reducing the image blurring.
Figure 6 demonstrates the results of the algorithm in modifying an image that was artificially and exaggeratedly fogged.


Figure 6. (a) Artificially and exaggeratedly fogged image; (b) corrected version of fogged image.

### 2.3. Determining the Exact Location of the License Plate

After removing the adverse environmental effects, according to the following two principles, the exact location of the license plate must be determined: First, due to the difference in the colors around the license plate and its background, by using edge detection on the black-white image the edge of the license plate frame will appear as an edge and a closed path shape. Secondly, according to each country's standards, the length-to-width ratio of the license plate will be a fixed value.

According to the above content, all the edges on the image are detected. Detected edges become bolder to remove any interruptions in the closed paths. Then, by defining the closed path edges in the image as objects, the one with the standard license plate's length-to-width ratio is selected as the main object (the license plate frame). Detecting the main object's position from the initial image determines the license plate frame. Figure 7 shows all the above steps on a sample image.


Figure 7. Detecting the license plate frame: (a) the initial image; (b) edge detection for black-white mode of initial image; (c) filling of closed-path detected edges; (d) the main objects in the image and finding the object corresponding to the standard license plate; (e) detecting the main object's position from the initial image determined the license plate frame.

### 2.4. Determining the Segments inside the Plate

After cutting the image of the license plate from the original image, by applying rotation if needed, removing the unessential edges of the license plate, and turning it into the black and white mode, the segments of recognizable license plate characters are separated [25-29].

In the image shown in Figure 8, from the left side, the index of the first column has at least one white pixel, labeled as the start index of the first segment. Additionally, the index of the first column without a white pixel is labeled the final index of the first segment. This process is repeated for the whole of the plate to determine all its character segments. Figure 9 shows the cut segments separately.

## 45E933 10

Figure 8. License plate image prepared for extracting its character segments.

## 4

Figure 9. Cut segments of the license plate-black and white mode image.

### 2.5. Recognizing the Segments Using Hopfield's Neural Network

Determined character segments in the previous section should be recognized using standard character patterns. Diverse methods of pattern recognition include matching pixel by pixel [26], the k-nearest neighbor, and Bayesian, and various neural networks can be used [30]. The Hopfield neural network is known as the most common method for detecting patterns with binary features. Since the extracted black and white segment images have binary values (zero for black and one for white), they can be recognized using this neural network [31].

The main idea of the Hopfield neural network is based on state variables. If the new position of a system depends on its previous one, it can be written in terms of state variables in the form of the following equation [32-35]:

$$
\begin{equation*}
x(t+1)=f(x(t)) \tag{10}
\end{equation*}
$$

The sequence above continues until its energy is exhausted, and then remains in a balanced state. The system energy should decrease to reach this state. For this purpose, as shown in Figure 10, the Hopfield neural network is designed such that, firstly, the new position of the system is dependent on its previous one, and second, its energy equation decreases.


Figure 10. A typical Hopfield neural network; its position depends on the previous time's position and its energy is decreasing.

The energy function and its gradient for the system shown in Figure 10 are defined as:

$$
\begin{align*}
E & =-\frac{1}{2} \sum_{i} \sum_{j} \omega_{i j} x_{i} x_{j}-\sum_{i} b_{i} x_{i} \\
\frac{d E}{d t} & =\sum_{i} \frac{d E}{d x_{i}} \frac{d x_{i}}{d t} \\
& =\sum_{i} \frac{d E}{d x_{i}} \frac{d x_{i}}{d n e t_{i}} \frac{\text { dnet }_{i}}{d t} \\
& =\sum_{i}\left(-\frac{d_{n e t}^{i}}{d t}\right) \frac{d x_{i}}{d n e t_{i}} \frac{\text { dnet }_{i}}{d t}  \tag{11}\\
& =-\sum_{i}\left(\frac{\text { dnet }_{i}}{d t}\right)^{2} \frac{d x_{i}}{\text { dnet }_{i}} \\
& =-\sum_{i}\left(\frac{\text { dnet }_{i}}{d t}\right)^{2} \varphi^{\prime}\left(\text { net }_{i}\right)
\end{align*}
$$

where $\varphi$ is the function of the neurons in the network and net ${ }_{i}$ is the output of the $i$ th neuron. According to Equation (11), the energy gradient of the system decreases if the derivative of the neuron function is positive. Choosing sign function as a function of the neuron can meet this condition. When Hopfield's neural network is used as an image classifier, the two-dimensional images should be mapped to one-dimensional mode, and the black pixel values set to -1 (while the value of the white pixels are +1 ).

The segments which are in black and white format are resized to a standard size. The matrix of each segment, which is two-dimensional, is transformed into a one-column vector. The black pixels marked with 0 in this matrix are changed to -1 and then applied to the Hopfield neural network. On the other hand, in this neural network, the main characters in standard size and 0 values corresponding to black pixels, which are replaced by -1 , are defined as balanced points. The Hopfield neural network moves the input matrix to the nearest balanced state. In other words, the closest standard character similar to the target segment is recognized. Figure 11 shows the noise (caused by mud) on the license plate. The designed Hopfield's neural network classified this noised license plate's character segments without error.


Figure 11. Detecting procedures of noisy license plate characters using Hopfield's neural network: (a) initial image; (b) detected license plate box; (c) the cut frame of the license plate; (d) black and white mode of cut frame; (e) cut noisy segments.

To determine the accuracy of Hopfield's neural network in determining numbers and letters, according to Figure 12, a set of different car license plates was considered. The graphics on these plates can play the role of noise. Hopfield's neural network classified the 253 characters on the license plates of this collection (after image processing and segmentation). Among the 253 test characters, only 6 characters were recognized wrongly, showing an accuracy of $97.6 \%$. Using a computer with Core(TM) i7-2640 CPU and 8 GB RAM, the time spent on determining the characters of each license plate is about 0.08 s .


Figure 12. A set of selected license plates from different standards to determine the CCR of the proposed algorithm.

In addition to the 253 main characters (numbers or letters), there are also 26 special characters (such as a dash or a combination of numbers and letters). According to the different standards in these plates, all special characters were considered as a unit pattern. In addition to correctly recognizing the main characters, Hopfield's neural network could also classify special characters in this unit pattern. The updated classification rate considering the special characters was $97.1 \%$.

The agile training capability of the Hopfield neural network has made it appropriate for application to plates with different standards, while for other neural networks, such as convolutional neural networks, a huge set of training data must be collected for each standard of the plates. Moreover, the accuracy of the proposed algorithm was higher than a number of similar ones developed on SOTA $[15,16]$. On the other hand, the time spent on recognizing the characters of each license plate was almost equal to the time spent on recognizing only one character in methods based on convolutional neural networks [15-18].

## 3. Conclusions

A new license plate pattern recognition system has been presented that is robust against adverse environmental effects such as fog or mud. Unlike previous studies that only considered a certain standard of license plates, this work evaluated all objects irrespective of their types. However, the selection of objects depends on their positions, and if the recognition of the license plate characters is unsuccessful another object enters the recognition process. Additionally, the paper has addressed the challenges posed by the presence of fog and smoke in the image by removing the matter of the image before initiating the license plate recognition process. More importantly, the use of Hopfield's neural network for license plate recognition, instead of the conventional method of pixel-to-pixel comparison of image segments with standard characters, has significantly reduced the execution time and increased the accuracy. The results pinpoint the efficacy of this approach. The neural network has shown capability in removing the noise on the license plate, making it a reliable tool for license plate recognition. The findings of this study contribute to the field of automated surveillance systems by providing an effective and efficient method for license plate recognition.

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## Article

# Application of Evolutionary Optimization Techniques in Reverse Engineering of Helical Gears: An Applied Study 

Vahid Pourmostaghimi ${ }^{1}$ (D) Farshad Heidari ${ }^{1}$, Saman Khalilpourazary ${ }^{2}$ and Mohammad Reza Chalak Qazani ${ }^{3,4, * *(D)}$<br>1 Department of Manufacturing and Production Engineering, Faculty of Mechanical Engineering, University of Tabriz, Tabriz 51666-16471, Iran<br>2 Department of Renewable Energy, Faculty of Mechanical Engineering, Urmia University of Technology, Urmia 57166-93188, Iran<br>3 Faculty of Computing and Information Technology, Sohar University, Sohar 311, Oman<br>4 Institute for Intelligent Systems Research and Innovation, Deakin University, Geelong, VIC 3216, Australia<br>* Correspondence: m.r.chalakqazani@gmail.com or m.chalakqazani@deakin.edu.au; Tel.: +61-411-857009

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#### Abstract

Reverse engineering plays an important role in the manufacturing and automobile industries in designing complicated spare parts, reducing actual production time, and allowing for multiple redesign possibilities, including shape alterations, different materials, and changes to other significant parameters of the component. Using reverse engineering methodology, damaged gears can be identified and modeled meticulously. Influential parameters can be obtained in the shortest time. Because most of the time it is impossible to solve gear-related inverse equations mathematically, metaheuristic methods can be used to reverse-engineer gears. This paper presents a methodology based on measurement over balls and span measurement along with evolutionary optimization techniques to determine the geometry of a pure involute of a cylindrical helical gear. Advanced optimization techniques, i.e., Grey Wolf Optimization, Whale Optimization, Particle Swarm Optimization, and Genetic Algorithm, were applied for the considered reverse engineering case, and the effectiveness and accuracy of the proposed algorithms were compared. Confirmatory calculations and experiments reveal the remarkable efficiency of Grey Wolf Optimization and Particle Swarm Optimization techniques in the reverse engineering of helical gears compared to other techniques and in obtaining influential gear design parameters.


Keywords: reverse engineering; helical gear; evolutionary optimization techniques; grey wolf optimization; whale optimization; particle swarm optimization; genetic algorithm

## 1. Introduction

Widely exploited in almost all engineering fields, reverse engineering (RE) is an approach that consists of digitizing a real component to create a numerical or virtual model [1]. This numerical or virtual model can also be emulated in presenting a Digital Twin. The Digital Twin, which started to be used recently, is a new concept for describing a new wave in modeling and simulation. While the use of simulation tools was previously restricted to design stages, nowadays different types of simulation tools are used in testing, validation, or optimization that can be referred to as Digital Twin [2]. This technology is an emerging concept that has become the center of attention in industry, particularly in the manufacturing industry. The Digital Twin is best described as the effortless integration of data between a physical and virtual machine in either direction [3]. Digital Twin has been widely used in different industries, especially manufacturing, to monitor performance, optimize progress, simulate results, and predict potential errors. This technique also plays various roles within the whole product lifecycle from design, manufacturing, delivery, use, and end of life. With the growing demands of individualized products and implementation of Industry 4.0, Digital Twin can provide an effective solution for future product design, development, and innovation [4]. Generally, RE can be defined as a principled process
that extracts design information from a product. Hence, components in various industries, including engineering, medical, defense systems, software, consumer electronics, etc., can be reverse-engineered comprehensively. Chikofsky and Cross introduced RE as "the process which analyzes a certain system to identify the systems" parts and their relationships and to create alternatives of the system in another form or at a higher level of abstraction [5]. In addition, the definition of RE in mechanical design is to "commence the redesign process where a product is observed, disassembled, analyzed, tested, and documented in terms of its functionality, form, physical principles, manufacturability, and assemble-ability.

In recent years, traditional approaches to designing, manufacturing, and constructing mechanical components and systems have been abolished. In other words, state-of-the-art methods in design and manufacturing have provided some outstanding benefits, including production preparation time decline, outstanding precision, etc., that completely change the quality of the product [6]. Nowadays, the implementation of RE is in manufacturing three-dimensional digital models of different mechanically damaged or broken components [7]. Lippmann et al. used RE for the verification of physical designs in nanoscale technologies [8]. RE is used in various fields, especially mechanics, such as composites, aerospace industries, energy plants, turbines, internal combustion components, etc.

One of the most prevalent applications of RE is redesigning and manufacturing gears with a high-precision involute widely used in the automobile industry. Modeling machine components, especially gears, with parametric or non-parametric RE methods is becoming more popular before producing samples using rapid prototype processes [9,10]. Numerous studies have been conducted on the redesign, modeling, and manufacturing of various types of gears through RE. Shamekhi et al. used RE to redesign and optimize the gear ratio and corresponding teeth number of the gears of an automatic transmission system. They proposed an accurate and efficient model to find optimum design parameters [11]. Verim et al. redesigned a damaged motor cam gear by the 3D scanning method. They evaluated the geometric values of the damaged model and the prototyped one. They reported acceptable deviations between the RE model and the actual component [12]. Dubravcik et al. made a 3D model of a damaged gear wheel using scanning and rapid prototyping. They compared and identified the new gear wheel's proportions and geometry with the original point cloud of the damaged gear. The results revealed good compliance between the original and RE-made gear [9]. Palkahas worked on the reconstruction of gears using RE and 3D printing methods. He reported that the final component's quality is highly related to the use of appropriate methods, software, hardware, and material [6]. Baehr et al. used machine learning for structural characterization in the RE domain [13].

The design of gears has a rather complex procedure involving numerous design parameters and factors. Finding an optimal or near-optimal solution for this complexity turns it into a challenging demand. There are plenty of studies in which different optimization algorithms, such as the Genetic Algorithm (GA) [14], Simulated Annealing (SA) [15], Response Surface Methodology (RSM) [16], Particle Swarm Optimization (PSO) [17], Grey Wolf Optimization (GWO), Whale Optimization Algorithm (WOA) [18], etc., have been implemented for the design optimization of gears. Xia et al. have used PSO combined with GA to optimize the operational quality of power shift transmission. The simulation and test results showed that the proposed control strategy effectively avoids the power cycle when shifting and consequently improves shift quality [19]. Artony has proposed a deterministic approach to solve simulation-based, multi-objective gear design optimization problems in the presence of general nonlinear constraints on the design variables. He has declared that, although developed with gear optimization in mind, it is broader in scope and should be tested in other computer-assisted engineering optimization problems [20]. Rai et al. used a Real-Coded Genetic Algorithm (RCGA) to achieve optimal helical gear design. They minimized the volume of a helical gear by increasing the profile shift coefficient as a design factor and other parameters such as the module, face width, and number of teeth using RCGA [21]. Mendi et al. investigated the difference between the dimensional optimization of motion and force-transmitting components of a gearbox performed by the
genetic algorithm and analytical method. The results illustrated that GA is a considerably better and more reliable method to obtain minimum gear volume [22].

Because of the drawbacks of the commonly performed research in the field of RE of gears, there is an urgent demand to carry out an extensive investigation to determine accurate gear design parameters using relevant measurements such as over-ball measurements. For instance, in the reviewed case study, the matter of RE for a single gear has not been investigated. In the studies dealing with this issue, parameters such as outside diameter and root diameter are considered known input parameters in the RE process. These practices pose an excessive error in the obtained design parameters. On the other hand, current industrially applicable gear design software requires initial design parameters such as the pressure angle and module to initiate the RE process. This problem adversely affects the accuracy of calculated design parameters. In industry, to replace a damaged gear wheel, it is commonplace to get its CAD model using scanning techniques. Many factors affect 3D scanning processes. One of them is the reflective ability of the parts' surfaces. Such inaccuracies and incompetence are mainly of the tooth systems because of uncontrolled laser ray scattering and incorrect scanning [10]. Therefore, a systematic methodology is needed to solve the RE problem of gears more comprehensively. However, the major requirement that has to be satisfied for successful gear RE is the modeling of the system at hand as accurately as possible [23].

Due to the complexities governing the calculations of input gear parameters, a new methodology was proposed in the current study to reverse-engineer gears using evolutionary optimization techniques. A number of commonly used metaheuristic optimization methods with successful applications reported in engineering optimization problems, such as GWO, WOA, PSO, and GA, were applied to solve the introduced RE problem. The stability, convergence speed, and accuracy of the mentioned algorithms were compared and evaluated. The paper has the following sections. A brief theoretical background of gear calculations is presented in Section 2. The proposed methodology and optimization techniques are explained in Section 3. Section 4 describes the remarkable results of the study. The conclusion is presented in Section 5.

## 2. Theoretical Background of Gear Calculations

The accurate design of gears requires considerable effort in calculating the relevant gear geometry to fulfil the essential requirements of operational characteristics such as load capacity, reliability, and gear size, which directly affect the quality of power transmission. The independent variables that define a gear are the number of teeth $z$, normal module $m_{n}$, normal pressure angle $\alpha_{n}$, addendum modification $x$, and helix angle $\beta$. Knowing the parameters mentioned, the involute form of a gear flank can be determined. The other dependent variables, including outside diameter, root diameter, and face width, are usually determined based on transmission load, which can be manipulated by the designer through design limitations and requirements [24]. The parameters that directly relate to the gear flank involute form and the parameters that do not have such a relationship are listed in Table 1.

Table 1. Division of the geometrical characteristics of gears.

|  | Input Parameters | Output Parameters |
| :---: | :--- | :--- |
| Direct relation with <br> involute form | Module $\left(m_{n}\right)$ <br> Normal pressure angle $\left(\alpha_{n}\right)$ | Span measurement $\left(s_{m}\right)$ <br> Helix angle $(\beta)$ |
| Addendum modification $(x)$ | Over-ball measurement $\left(d_{m}\right)$ <br> Chordal thickness measurement $\left(s_{j}\right)$ |  |
| No relation with <br> involute form | Face width | Outside diameter <br> Root diameter |

One of the most important parameters that greatly impacts the gear design characteristics, load capacity of the gear, and backlash of engaged gears is tooth thickness. Normal tooth thickness can be calculated as follows:

$$
\begin{equation*}
T T_{\mathrm{n}}=\frac{m_{\mathrm{n}} \pi}{2}+2 \times m_{\mathrm{n}} \times x \times \tan \left(\alpha_{\mathrm{n}}\right) \tag{1}
\end{equation*}
$$

There are direct and indirect methods for measuring tooth thickness, including chordal thickness measurement $\left(s_{j}\right)$, span measurement $\left(s_{m}\right)$, and over-ball measurement $\left(d_{m}\right)$. Among these methods, span and over-ball measurements are more industrially applicable because of their higher measurement precision and the comparatively lower equipment involved in the measurement process [25]. A schematic view of span and ball measurements is illustrated in Figure 1.


Figure 1. A schematic view of span measurement and over-balls measurement for a gear.
The value for $s_{m}$ based on the aforementioned input parameters can be calculated as follows [24]:

$$
\begin{equation*}
s_{m}=m_{n} \times\left[\left(z_{k}-0.5\right) \times \pi \times \cos \left(\alpha_{n}\right)+z \times \operatorname{inv}\left(\alpha_{n}\right) \times \cos \left(\alpha_{n}\right)+2 \times x \times \sin \left(\alpha_{n}\right)\right] \tag{2}
\end{equation*}
$$

in which $z_{k}$ is the number of spanned teeth. The function $\operatorname{inv}\left(\alpha_{\mathrm{n}}\right)$ is calculated according to Equation (3).

$$
\begin{equation*}
\operatorname{inv}\left(\alpha_{n}\right)=\tan \left(\alpha_{n}\right)-\left(\alpha_{n}\right) \tag{3}
\end{equation*}
$$

The parameter $d_{m}$ is obtained as follows:

$$
d_{m}=\left\{\begin{array}{cc}
d_{k}+\frac{m_{n} \times z \times \cos \left(\alpha_{n}\right)}{\cos (\varphi) \times \cos (\beta)} & z=\text { even }  \tag{4}\\
d_{k}+\frac{m_{n} \times z \times \cos \left(\alpha_{n}\right)}{\cos (\varphi) \times \cos (\beta)} \times \cos (90 / z) & z=\text { odd }
\end{array}\right.
$$

in which $d_{k}$ is the diameter of the measuring ball, and $\varphi$ is the pressure angle at the center of the measuring ball, which can be calculated according to Equation (5) [24].

$$
\begin{equation*}
\operatorname{inv}(\varphi)=\frac{d_{k}}{m_{n} \times z \times \cos \left(\alpha_{n}\right)}-\frac{\pi}{2 \times z}+\operatorname{inv}\left(\alpha_{\mathrm{t}}\right)+\frac{2 \times x \times \tan \alpha_{n}}{z} \tag{5}
\end{equation*}
$$

In Equation (5), $\alpha_{\mathrm{t}}$ is the transverse pressure angle and is defined as below:

$$
\begin{equation*}
\alpha_{t}=\tan ^{-1}\left(\frac{\tan \left(\alpha_{n}\right)}{\cos (\beta)}\right) \tag{6}
\end{equation*}
$$

Consequently, the value of $\operatorname{inv}\left(\alpha_{t}\right)$ is calculated as Equation (7) [24]:

$$
\begin{equation*}
\operatorname{inv}\left(\alpha_{t}\right)=\tan \left(\alpha_{t}\right)-\left(\alpha_{t}\right) \tag{7}
\end{equation*}
$$

## 3. Materials and Methods

The design of gears has become more of a necessity than ever before. In particular, gear design optimization and transmission accuracy have attracted more attention [26]. One of the most frequently used techniques in gear design is RE methods. These methods are necessary for remodeling or measuring damaged or non-damaged gears. RE also enables the design of complex mechanical components, effectively reducing production time and relevant prototyping costs. Using RE methods, damaged gear components can be calculated and modeled in a short time, and obtained models can be produced economically in a comparably short period [26].

### 3.1. Reverse Engineering of Cylindrical Helical Gear

When designing a new set of gears, there are many alternatives in the gear parameters selection. Therefore, the designer can freely choose among the various options. However, RE imposes a considerable challenge to the gear designer since the main problem is characterizing a previously manufactured gear whose geometry is unknown.

With a valid reference for the design, the designer must extract the exact geometry of the redesigned gear. For the completion of the RE of a cylindrical helical gear, it is obligatory to determine the number of teeth $z$, normal module $m_{n}$, normal pressure angle $\alpha_{n}$, addendum modification $x$, and helix angle $\beta$. According to Equations (2)-(7), identifying the primary gear input parameters based on measurements obtained from $s_{m}$ and $d_{m}$ is impossible to solve these equations. This arises from the complex calculations governing the relations between gear input parameters and measured output characteristics. So, there is a great need to use novel methods in the RE of gears to obtain influential input parameters. Because of the inability of existing computational and mathematical methods to analytically solve these equations, it seems that applying modern evolutionary optimization algorithms can be a good choice to obtain gear design parameters in the RE process.

Inspired by natural processes, these techniques can optimize a function or find the answer to a complex equation in a logical period and with relatively high precision.

### 3.2. Problem Description and Solving Method

### 3.2.1. Objective Function

The general constrained problem considered in the present research was expressed as:

$$
\begin{gather*}
\text { minimize } E\left(m_{n}, \alpha_{n}, \beta, x\right) \\
\text { subject to } m_{n_{\min }} \leq m_{n} \leq m_{n_{\max }} \\
\alpha_{n_{\min }} \leq \alpha_{n} \leq \alpha_{n_{\max }}  \tag{8}\\
\beta_{\min } \leq \beta \leq \beta_{\max } \\
x_{n_{\min }} \leq \alpha_{n} \leq x_{n_{\max }}
\end{gather*}
$$

The objective function $E\left(x_{1}, x_{2}, x_{3}, x_{4}\right)$ to find input design parameters was set as follows:

$$
\begin{equation*}
E\left(m_{n}, \alpha_{n}, \beta, x\right)=\sum_{i=1}^{n} \sqrt{\left(\text { Formula obtained value }{ }^{i}-\text { Measured value }\right)^{i}} \tag{9}
\end{equation*}
$$

Since, in the present study, three values for $d_{m}$ and one value for $s_{m}$ are measured, Equation (9) can be rewritten as:

$$
\begin{equation*}
E\left(m_{n}, \alpha_{n}, \beta, x\right)=\sqrt{\left(s_{m}-S_{\text {measured }}\right)^{2}}+\sum_{i=1}^{3} \sqrt{\left(d_{m}^{i}-D_{\text {measured }}^{i}\right)^{2}} \tag{10}
\end{equation*}
$$

in which $S_{\text {measured }}$ and $D_{\text {measured }}$ are practically measured span and over-balls measurements, respectively. According to Equation (8), it can be concluded that the defined objective function is the resulting error of the selected design parameters and measured values. Therefore, the minimum possible value for $E\left(m_{n}, \alpha_{n}, \beta, x\right)$ can lead to the exact input gear design parameters. In this way, the present investigation strived to substitute an unsolvable series of complicated equations with an optimization problem to discover the answers by exploiting evolutionary optimization algorithms.

### 3.2.2. Methodology

Figure 2 presents a schematic framework that integrates the described RE technique with evolutionary optimization algorithms. The raw measured data are fed into the optimization unit to minimize the defined objective function. After the termination of the optimization algorithm, based on measured values and formulas, the optimum results are gained.


Figure 2. Flowchart of the proposed RE of gear framework.
Several termination conditions can be used to terminate the algorithm and obtain calculated optimum results. However, two are more frequently used for the problems under consideration, including the maximum number of function evaluations and minimum specific error. The former is suitable for computationally demanding problems when only a limited number of function evaluations is allowed. At the same time, the latter specifies that the optimization is terminated when the minimum desired error is reached [20].

Relatively high accuracy is required to determine the parameters in this research. Therefore, the algorithm is terminated in two modes, after a definite number of function evaluations and after reaching a predefined error. The minimum value for the cost function, $\varepsilon$, reflects the resolution of searching for optimum answers in the space of all possible
answers. The proposed algorithmic framework is illustrated by the pseudocode in Figure 3, which summarizes the methodology discussed.

```
Input:
    - Gear formulas for \(s_{\mathrm{m}}\) and \(d_{\mathrm{m}}\)
    - Lower and upper bounds for \(m_{\mathrm{n}}, \alpha_{\mathrm{n}}, \beta, x\)
    - The minimum value for cost function \(E\left(m_{\mathrm{n}}, \alpha_{\mathrm{n}}, \beta, x\right)\) to terminate the algorithm, \(\varepsilon\)
    - Selecting proper values for \(z_{\mathrm{k}}\) and \(d_{\mathrm{ki}}\) for \(\mathrm{i}=1: 3\)
    - Measuring \(s_{\mathrm{m}}\) and \(d_{\mathrm{mi}}\) for \(i=1: 3\)
1. Set \(\mathrm{i}=0\)
2. Create objective function \(E\left(m_{\mathrm{n}}, \alpha_{\mathrm{n}}, \beta, x\right)\)
3. do while \(E>\varepsilon\)
    3.1. Set \(\mathrm{F}=\mathrm{i}+1\)
    3.2. Generate initial population ( \(m_{\mathrm{ni}}, \alpha_{\mathrm{ni}}, \beta_{\mathrm{i}}, x_{\mathrm{i}}\) )
    3.3. Calculate \(s_{\mathrm{m}}\) and \(d_{\mathrm{m}}\) for members of population
    3.4. Calculate \(E\left(m_{\mathrm{n}}, \alpha_{\mathrm{n}}, \beta, x\right)\) according to Equation 9 .
    3.5. if \(E<\varepsilon\) then
        exit (leave the loop)
        else if
            Apply algorithm operators to create next population
    End
4. Return best member as gear design parameters ( \(m_{\mathrm{nopt}}, \alpha_{\mathrm{nopt}}, \beta_{\mathrm{opt}} x_{\mathrm{opt}}\) )
```

Figure 3. Pseudocode of the proposed algorithmic framework.

### 3.2.3. Evolutionary Optimization Algorithms

Although they have been used widely in various design problems, conventional optimization methods impose some serious difficulties on the matter; accordingly, to have a slow convergence speed, these methods may be stuck in the local optimum point. Furthermore, when the objective function or constraints cannot be stated functionally explicit in the input parameters, it is hard to utilize these methods to solve the given problem. Consequently, some novel methods, quite often inspired by nature, such as GA, PSO, SA, etc., are applied to solve complex optimization problems [22]. In this research, an investigation was carried out to find the design parameters of a gear based on performed measurements. The proposed algorithms were GWO, WOA, PSO, and GA, which were used under similar conditions and objective functions. The population size was set to 40 for all of the algorithms. The primary population was created as uniformly as possible within the search space. The results of the methods above were compared, and their competencies to solve the problem were evaluated.

Metaheuristic methods have been proved to be a powerful technique to solve complex optimization problems. The majority of them consider a set of parameters that must be tuned before the triggering of the algorithm or during the process and are based on multiple runs of the metaheuristic algorithm [27]. During each iteration, significantly worse configurations are removed, and new configurations are formed through crossover and mutation. Each of the presented methods has its own features as well as advantages and disadvantages [28].

A good setup of these parameter values can result in a better and effective application of the optimization process. This job can be carried out by trial-and-error methodology [29] or by intelligent tuning methods [27]. In this regard, various tuning methods have been presented such as CRS [28], F-Race [30], Revac [31], and ParamILS [32]. Control parameters for selected optimization algorithms in this research were selected based on a trial-and-error approach.

Grey Wolf Optimization Algorithm: GWO is a metaheuristic optimization algorithm that imitates the intelligent foraging behavior of grey wolves when hunting prey. Strictly following their social hierarchy, grey wolves are divided into 4 groups from top to bottom. The top group which makes decisions during hunting is called $\alpha$. The next level belongs to $\beta$ wolves who help as deputy chiefs in the group. The $\alpha$ wolves are replaced by $\beta$ ones after dying or becoming inefficient as leaders. In the third level are $\delta$ wolves, which act as hunters and scouts of the group. The last level in the established hierarchy is occupied by the weakest members, titled $\omega$ wolves. When hunting, all types of wolves are organized according to decisions made by $\alpha$ wolves to identify, follow the prey, encircle it, and finally attack it. In this optimization algorithm, $\alpha, \beta$, and $\delta$ indicate the members with fitness from high to low, in that order [33].

The simulation of wolves' behavior when hunting includes encircling and catching the prey. The mathematical description of encircling the prey can be expressed as follows:

$$
\begin{gather*}
\vec{D}=\left|\vec{C} \cdot \vec{X}_{p}(t)-\vec{X}(t)\right|  \tag{11}\\
\vec{X}(t+1)=\vec{X}_{p}(t)-\vec{A} \cdot \vec{D} \tag{12}
\end{gather*}
$$

in which $\vec{A}$ and $\vec{C}$ act as coefficient vectors. $\vec{X}_{p}$ determines the position of the prey or global answer. The vector $\vec{X}$ describes the current position of the wolves; the new position of $\alpha, \beta$, and $\delta$ wolves in the $t$ th iteration would be determined by $\vec{D}$. The values for $\vec{A}$ and $\vec{C}$ can be calculated as below:

$$
\begin{gather*}
\vec{A}=2 \cdot \vec{a} \cdot \vec{r}_{1}-\vec{a}  \tag{13}\\
\vec{C}=2 \cdot \vec{r}_{2} \tag{14}
\end{gather*}
$$

where $\vec{r}_{1}$ and $\vec{r}_{2}$ are random vectors varying in the range of $[0,1]$, and $\vec{a}$ is the convergence vector value changing linearly in $[0,2]$.

In the next step, the hunting process starts with $\alpha$ wolves as the leaders and $\delta$ wolves as the hunting contributors. Mathematically simulated hunting assumes that $\alpha, \beta$, and $\delta$ individuals have better information about the location of prey. Therefore, all individuals update their positions according to the below equations:

$$
\begin{align*}
& \vec{X}_{1}=\vec{X}_{\alpha}-A_{1} \cdot\left|\vec{C}_{1} \cdot \vec{X}_{\alpha}-\vec{X}\right|  \tag{15}\\
& \vec{X}_{2}=\vec{X}_{\beta}-A_{2} \cdot\left|\vec{C}_{2} \cdot \vec{X}_{\beta}-\vec{X}\right|  \tag{16}\\
& \vec{X}_{3}=\vec{X}_{\delta}-A_{3} \cdot\left|\vec{C}_{3} \cdot \vec{X}_{\delta}-\vec{X}\right|  \tag{17}\\
& \vec{X}(t+1)=\frac{\vec{X}_{1}+\vec{X}_{2}+\vec{X}_{3}}{3} \tag{18}
\end{align*}
$$

At the ones where the prey stops, the wolves attack. Mathematically, this happens when $|\vec{A}| \leq 1$. Unless the attack is delayed until finding a better position or answer [34], the schematic position update of members in the GWO algorithm is illustrated in Figure 4. The population size is set to 40 .


Figure 4. Position update in GWO.
Whale Optimization Algorithm: Introduced by Mirjalili et al., WOA is a swarm optimization method inspired by the hunting behavior of humpback whales. The hunting process involves searching, encircling the prey, and finally attacking to catch it [35]. The relatively easy implementation of the algorithm and its configuration with fewer parameters are the advantages of this optimization technique. The prey is the answer to the optimization problem that must be determined during the WOA process [36]. The hunting of humpback whales is illustrated in Figure 5. The movement of randomly created whales toward the prey or leader whale is shown in Figure 5a. The encircling process of the prey by the whales and their spiral movement while emitting bubbles to surround the prey is illustrated in Figure 5b,c, respectively [37].


Figure 5. Hunting process of humpback whales: (a) movement of whales in a random direction or toward leader; (b) encircling process after finding the prey; (c) reaching the prey in spiral route.

First, a group of whales or members is created randomly to search for the prey or answer in the search space. This stage can be expressed mathematically as follows:

$$
\begin{equation*}
\vec{X}(t+1)=\vec{X}_{\text {rand }}(t)-\vec{A} \cdot\left|\vec{C} \cdot \vec{X}_{\text {rand }}(t)-\vec{X}(t)\right| \tag{19}
\end{equation*}
$$

in which $\vec{A}$ and $\vec{C}$ act as coefficient vectors. $\vec{X}_{\text {rand }}(t)$ determines the position of whales in tth iteration. The values for $\vec{A}$ and $\vec{C}$ can be obtained as follows:

$$
\begin{gather*}
\vec{A}=2 \cdot \vec{a} \cdot \vec{r}-\vec{a}  \tag{20}\\
\vec{C}=2 \cdot \vec{r}  \tag{21}\\
a=2-\frac{2 t}{t_{\max }} \tag{22}
\end{gather*}
$$

In these equations, $\vec{r}$ is a vector selected randomly in $[0,1]$. The maximum number of permissible iterations is $t_{\text {max }}$, and $\vec{a}$ is the convergence vector value which decreases linearly from 2 to 0 by increasing the number of iterations.

After having found the prey, the whales move closer to the prey. The mathematical formula of this shrinkage mechanism can be expressed as follows:

$$
\begin{equation*}
\vec{X}(t+1)=\vec{X}_{g-b e s t}-\vec{A} \cdot\left|\vec{C} \cdot \vec{X}_{g-b e s t}-\vec{X}(t)\right| \tag{23}
\end{equation*}
$$

It must be added that if $|\vec{A}| \leq 1, \vec{X}_{g-b e s t}$ determines the member with the best fitness value. Therefore, $\vec{X}_{g-b e s t}$ will update the position of members automatically to encircle the prey. In the final step, the whales catch their prey, start to move in a spiral route, and simultaneously narrow the siege [38]. This process is illustrated in Figure 6.


Figure 6. Search mechanism and whales' movement in WOA: (a) shrinking encircling process; (b) position update in the spiral way.

Figure 6a demonstrates a probable number of positions such as $(X, Y)$, which can be achieved by the current global optimum $\left(X^{*}, Y^{*}\right)$ in any determined way. The spiral way to narrow the encirclement is shown in Figure 6b [37]. This behavior can be mathematically expressed as below:

$$
\begin{equation*}
\vec{X}(t+1)=\left|\vec{X}_{g-b e s t}-\vec{X}(t)\right| \cdot e^{b l} \cdot \cos (2 \pi l)+\vec{X}_{g-b e s t} \tag{24}
\end{equation*}
$$

in which $b$ determines the shape of the spiral route, and $l$ is a random value in the range of [0, 1] [39]. The population size is set to 40 .

Particle Swarm Optimization: Nature-inspired metaheuristic optimization algorithms have recently been extensively utilized to optimize various manufacturing problems [40].

PSO has gained a special place among these algorithms because of its simplicity in programming and solving relatively complex functions [41].

Introduced by Eberhart and Kennedy in 1995 [42], PSO is a population-based optimization algorithm which comprises several particles representing a probable solution to the defined optimization problem. The goal of these particles in the algorithm is to find the answer by improving their positions [43].

Particles have a fitness factor, which determines the ability of each particle to solve the optimization problem, and a velocity factor which affects the movement direction of a particle in the course of iterations and can be mathematically obtained as follows:

$$
\begin{equation*}
V_{i d}^{(t+1)}=w V_{i d}^{(t)}+c_{1} \operatorname{rand}_{1}\left(p_{\text {best } i d}^{(t)}-X_{i d}^{(t)}\right)+c_{2} \operatorname{rand}_{2}\left(g_{\text {best } i d}^{(t)}-X_{i d}^{(t)}\right) \tag{25}
\end{equation*}
$$

where $X_{i d}^{(t)}$ and $V_{i d}^{(t)}$ stand for the position and velocity of particle $i$ in $d$ dimensional space, respectively. $p_{\text {best id }}^{(t)}$ and $g_{\text {best id }}^{(t)}$, respectively, give the best position of particle i and other particles in the population until generation $t$. Inertia weight factor $w$ adjusts the dynamic behavior of particles. Parameters rand1 and rand2 are random variables in the range of $[0,1], c 1$ is the determined cognitive factor, and c2 is the preferred social factor of each particle. The positions of particles are updated considering the calculated velocity:

$$
\begin{equation*}
X_{i d}^{(t+1)}=X_{i d}^{(t)}+V_{i d}^{(t-1)} \tag{26}
\end{equation*}
$$

in which $X_{i d}^{(t+1)}$ and $X_{i d}^{(t)}$ are a new position and previous position of particle i. The updating process of particles' positions is illustrated in Figure 7.


Figure 7. Updating the position of particles in PSO.
The optimization process terminates when the best answer is found, or the desired iteration is reached. The parameter configuration selected in the PSO implementation process in the research is shown in Table 2. The data were selected based on a trial-anderror method to reach the minimum possible error.

Table 2. PSO parameter configuration.

| Population size | 40 |
| :---: | :---: |
| Range of inertia weight | $0.4-0.8$ |
| Cognitive factor | 1.5 |
| Social factor | 1.5 |
| Stopping criteria | Minimum Specified Error |

Genetic Algorithm: As the most recognized subset of evolutionary algorithms, it mimics the biological evolution of organisms in nature. Having been utilized in various engineering problems, GA was initially developed and characterized by John Holland [44]. This naturally inspired optimization algorithm is structured on natural genetics and selection. It allows only the solutions to survive and produce successive generations with maximum fitness according to the defined optimization problem [45].

GA deals with a set of members or solutions in the total generated population. The solutions are ranked in the population based on their fitness values. The solutions placed at the top have more chance to participate in the reproduction of the next generations, and other solutions with lower fitness must be removed from the population, and new solutions can be generated. This happens by applying crossover and mutation operators to individuals. A sample of operators applied to a pair of parents to generate new offspring is illustrated in Figure 8.


Figure 8. Crossover and mutation operators in GA.
Parameter configurations of GA in this research are given in Table 3. These parameters were obtained through the trial-and-error method.

Table 3. Parameter configuration of GA.

| Population size | 40 |
| :---: | :---: |
| Length of chromosomes | 6 |
| Selection operator | Roulette wheel |
| Crossover operator | Single-point operator |
| Crossover probability | 0.7 |
| Mutation probability | 0.15 |
| Fitness parameter | Operation time |

### 3.2.4. Experimental Methodology

For evaluating the proposed RE methodology of a sample gear, a helical gear related to the 1st forward movement gear of a Massey Ferguson 399 tractor's gearbox was selected. The drawing specifications, including the input-independent design parameters of the selected transmission gear, are given in Table 4.

The output measurable parameters of the gear, including values for span size and over-two-balls measurement, are listed in Table 5. The values were validated using special gear software. The measurements over two balls and span size using appropriate balls, a caliper, and a disc micrometer are demonstrated in Figure 9. To determine the form of involute, three values for $d_{m}$ with different balls were obtained. The smallest one fell below the pitch diameter of the gear, the medium one was closer to the pitch diameter, and the bigger pin diameter was used to encompass the involute area of the gear over the pitch
diameter. Moreover, a span measurement $s_{m}$ was used to decrease the potential error of calculations and increase the precision of obtained results.

Table 4. Input-independent design parameters of the selected gear.

| Input Parameter | Values of the Selected Gear Part |
| :---: | :---: |
| $z$ | 42 |
| $m_{n}$ | 4.233 |
| $\alpha_{n}$ | 25 |
| $\beta$ | $17^{\circ} 59^{\prime} 58^{\prime \prime}$ |
| $x$ | +0.0863 |

Table 5. Measurements of the selected gear.

| Output Parameter | Condition of Measurement | Measured Value |
| :---: | :---: | :---: |
| $d_{m}^{1}$ | $d_{k}=6.5$ | 195.264 |
| $d_{m}^{2}$ | $d_{k}=7$ | 196.890 |
| $d_{m}^{3}$ | $d_{k}=7.5$ | 198.493 |
| $s_{m}$ | $k=7$ | 87.206 |



Figure 9. Measuring of over two balls and span size.
Since $\alpha_{n}$ is usually an integer value in the design process of gears, the value of this parameter is taken with no decimal while applying the optimization algorithm to find the valid answer in the search space. The rest of the parameters, including $m_{n}, \beta$, and $x$, are considered with three decimal places.

The optimization is performed in two scenarios. In the first mode, the algorithm terminates after reaching 10,000 function evaluations. The algorithms continue in the latter until the minimum desired error of 0.005 is reached. This error value guarantees that the calculated input design parameters would be obtained with the highest possible accuracy and the least deviation from actual input parameters.

All evolutionary optimization algorithms were triggered in MATLAB software. The minimum and maximum values for design parameters in search space are given in Table 6. These limits have been selected according to widespread applications observed in the power transmission mechanisms used in the automotive industry.

Table 6. The limit values of input design parameters in search space.

| Parameter | Minimum | Maximum |
| :---: | :---: | :---: |
| $m_{n}$ | 1 | 5 |
| $\alpha_{n}$ | 10 | 30 |
| $\beta$ | 0 | 30 |
| $x$ | -1 | 1 |

## 4. Results and Discussion

A computer equipped with Intel ${ }^{\circledR}$ core TM i7, 4.00 GHz CPU, and 16 GB RAM was selected to run algorithms in MATLAB ${ }^{\circledR}$ software, version 2020b. For optimization with the GA, function $g a$ in MATLAB was used. For optimization with the PSO algorithm, Particle Swarm Optimization MATLAB Toolbox Version 1.0.0.0 was selected. For optimization with the GWO algorithm, Grey Wolf Optimizer MATLAB Toolbox Version 1.0 was utilized, and for optimization with WOA, Whale Optimization Algorithm MATLAB Toolbox Version 1.0 was applied. The population size for all algorithms was equal to perform a fair comparison of the accuracy and convergence speed of the algorithms. The performance of all optimization algorithms was submitted after 20 runs, and the run with the best result was selected in comparison with the efficiency of the algorithms. The parameter configurations for GA and PSO were discussed in the corresponding sections. GWO and WOA do not need any specific parameters [13]. As mentioned before, the RE problem was approached with two different scenarios. Concerning the first one, the efficiency of the algorithms in terms of their convergence speed was assessed by selecting a definite number of function evaluations for each algorithm to run. In the second scenario, the minimum error was selected as the final goal of optimization. Overall, the fastest algorithm to find optimum input design parameters will be determined in this case.

Scenario I: After performing 10,000 function evaluations, the algorithm's performance in finding the best solution for the problem was evaluated. The results of the optimizations and the time elapsed for each algorithm to find its best solution are given in Table 7. While the PSO method presents the best solution, other algorithms offer different values for the problem. The best result for GWO was obtained after 8824 function evaluations within 549.8 s . This means that GWO is superior to other algorithms in terms of convergence speed and proximity to the best answer. The WOA, PSO, and GA algorithms followed in obtained accuracy with elapsed times of $736.4 \mathrm{~s}, 1122.3 \mathrm{~s}$, and 2364.1 s , respectively.

Table 7. Best solutions of the proposed algorithms for scenario I and relevant obtained input design parameters.

|  | GWO | WOA | PSO | GA |
| :---: | :---: | :---: | :---: | :---: |
| Normal module $\left(m_{n}\right)$ | 4.2333 | 4.0012 | 4.2122 | 5.2666 |
| Normal pressure angle $\left(\alpha_{n}\right)$ | $25^{\circ}$ | $25^{\circ}$ | $25^{\circ}$ | $21^{\circ}$ |
| Addendum modification $(x)$ | 0.0863 | -0.0128 | 0.0055 | 0.1992 |
| Helix angle $(\beta)$ | $17^{\circ} 59^{\prime} 58^{\prime \prime}$ | $14^{\circ} 36^{\prime} 32^{\prime \prime}$ | $17^{\circ} 02^{\prime} 11^{\prime \prime}$ | $13^{\circ} 52^{\prime} 20^{\prime \prime}$ |
| Function evaluations | 8824 | 10,000 | 10,000 | 10,000 |
| Elapsed time | 549.8 | 736.4 | 1122.3 | 2364.1 |

The values obtained by the algorithms given in Table 7 show that the error of PSO and WOA is $2.2 \%$ and $12.3 \%$ greater than the expected value of 0.003 . This value is $41.9 \%$ for GA, which indicates the weakness of this method compared to the other swarm-based optimization techniques utilized in the present research. The error was calculated based on Equation (10) for the input design parameters for each algorithm and measured values. The obtained design parameters for each algorithm are also given in Table 7.

In Table 8, the statistical results of the performance of each algorithm in all the runs are summarized. The table contains the average number of function evaluations, elapsed computation times, and the deviation of obtained results.

Table 8. Statistical results of the performance of each algorithm in all runs.

|  | GWO | WOA | PSO | GA |
| :---: | :---: | :---: | :---: | :---: |
| Function evaluations (ave.) | 9502 | 10,000 | 10,000 | 10,000 |
| Elapsed time (ave.) | 601.5 | 751.1 | 1206.7 | 2404 |
| Standard deviation | 0.1421 | 4.0255 | 0.9054 | 29.1512 |
| Function evaluations (ave.) | 9502 | 10,000 | 10,000 | 10,000 |

The values given in Table 8 reveal that the repeatability of the algorithms can differ drastically. While PSO and GWO demonstrate relatively stable repeatability, the other two algorithms have negligible values, attributed to their inadequate efficiency in the optimization process in the RE of gears.

Scenario II: In this case, all algorithms were expected to fulfill predefined error criteria. The aim was to compare the efficiency of the proposed algorithms in terms of their convergence speed in finding the best solution. The results of the optimizations, including the optimum design parameter, number of function evaluations, and time elapsed for each algorithm to find its best solutions, are given in Table 9. While PSO, GWO, and WOA reached the best solution with differences in the function evaluations and consequently the computation time, the best answer was not reached by all the runs for GA. The results show the superiority of swarm-based methods in optimizing complicated problems. GWO has the highest convergence speed and lowest number of function evaluations among the proposed swarm-based optimization methods. PSO and WOA are placed in the next levels with $127.4 \%$ and $228.7 \%$ more required time and a $26.2 \%$ and $178.1 \%$ higher number of function evaluations, respectively.

Table 9. Best solutions of the proposed algorithms for scenario II and relevant obtained input design parameters.

|  | GWO | WOA | PSO | GA |
| :---: | :---: | :---: | :---: | :---: |
| Normal module $\left(m_{n}\right)$ | 4.2333 | 4.2333 | 4.2333 | 4.0882 |
| Normal pressure angle $\left(\alpha_{n}\right)$ | $25^{\circ}$ | $25^{\circ}$ | $25^{\circ}$ | $25^{\circ}$ |
| Addendum modification $(x)$ | 0.0863 | 0.0863 | 0.0863 | 0.0012 |
| Helix angle $(\beta)$ | $17^{\circ} 59^{\prime} 58^{\prime \prime}$ | $17^{\circ} 59^{\prime} 58^{\prime \prime}$ | $17^{\circ} 59^{\prime} 58^{\prime \prime}$ | $14^{\circ} 12^{\prime} 10^{\prime \prime}$ |
| Function evaluations | 8824 | 24,541 | 11,142 | - |
| Elapsed time | 549.8 | 1807.1 | 1250.5 | - |

The statistical results of the performance of each algorithm in all the runs are given in Table 10. In the table, each algorithm's best, mean, and worst performance; average number of function evaluations; and average elapsed time to terminate the algorithms are given.

Table 10. Statistical results of the performance of each algorithm in all runs.

|  | GWO | WOA | PSO | GA |
| :---: | :---: | :---: | :---: | :---: |
| Function evaluations (ave.) | 9816 | 34,101 | 12,294 | - |
| Elapsed time (ave.) | 712.5 | 2511 | 1892.1 | - |
| Standard deviation | 0.1902 | 8.9778 | 0.2252 | - |

According to the data given in Table 10, the GWO algorithm with the lowest standard deviation is superior to the other considered algorithms in terms of stability and robustness. In all performed runs, GWO showed relatively uniform performance in reaching the best solution. Furthermore, based on Table 10, GWO, with quite a short elapsed time, is the fastest algorithm compared to the others.

The convergence graphs by the considered algorithms are shown in Figure 10. The curves are derived from the best performance of each algorithm during the optimization run.

All the methods except GA obtained the optimum input design parameters. GWO and PSO converge to the optimum values much more quickly than others without being stuck in the local minimum solution. However, the convergence speed of GWO is relatively higher than that of PSO. This method finds the optimum value in fewer function evaluations. The performance of WOA is relatively weak in this scenario. Being attached to a local minimum during the optimization progress requires more time and function evaluations for the WOA to reach the optimum value. Among the proposed optimization algorithm, GA has the lowest convergence performance. No optimum results were found in any of the runs for

GA. Therefore, it can be concluded that GA is attached to the local optimum solution and cannot be considered as an appropriate algorithm for the defined problem.


Figure 10. The algorithms' convergent curves reaching the minimum defined objective function for GWO, WOA, PSO, and GA.

To consider the average number of iterations in which the best results for each proposed algorithm are reached, boxplots for GWO, WOA, and PSO were generated and are shown in Figure 11. An illustrative comparison via the given boxplots demonstrates the superiority of GWO over the other two algorithms. Because of the poor performance of GA in finding optimum results, it has been omitted from the performance comparison.


Algorithm
Figure 11. Boxplots of the iterations for GWO, WOA and PSO.
A two-sample unpaired $t$-test was performed to test the hypothesis that the mean number of iterations for GWO is lower than the same for the WOA and PSO algorithms. $t$-tests are divided into two groups. The first group is an unpaired $t$-test, which can be used when the two groups under comparison are independent of each other, and the second is a paired $t$-test, which can be used when the two groups under comparison are dependent on each other. As the $t$-test is a parametric test, samples should meet certain preconditions, such as normality and equal variances. In other words, in order to reach a
statistical conclusion about a sample mean with a t-distribution, certain conditions must be satisfied: the two samples under comparison must be independently sampled from the same population, satisfying the conditions of normality and equal variance. For checking the normality and equality of variance, a Shapiro-Wilk test and Levene's test were utilized, respectively. According to the calculated Shapiro-Wilk test, the obtained $p$ value is 0.193 . Since the Shapiro-Wilk $p$ value is greater than 0.05 , it can be concluded that the samples satisfy the condition of normality. Moreover, the results of the Levene's test show that the condition of equal variance is also met.

The results for the $t$-test are given in Table 11. It is obvious from the test results that the zero hypothesis is rejected, and it can be concluded that there will be a meaningful difference between the effectiveness and convergence speed of the optimization algorithms in terms of the mean number of iterations, as given in Table 11.

Table 11. Two-sample t-test results for used algorithms based on iterations.

|  | $\boldsymbol{T}$-Test for GWO-WOA | $\boldsymbol{T}$-Test for GWO-PSO |
| :--- | :---: | :---: |
| $p$ value | $<0.0001$ | 0.0005 |
| Significant difference $(p<0.05)$ | YES | YES |
| $t$ value | 12.80 | 4.211 |
| df | 18 | 18 |
| Mean iterations | $10,957 / 28,800$ | $10,957 / 15,722$ |
| Difference between mean values | $17,843 \pm 1394$ | $4765 \pm 1132$ |
| $95 \%$ confidence interval | 14,913 to 20,773 | 2388 to 7143 |

According to the performed $t$-tests, the hypothesis claiming the higher effectiveness of the GWO over WOA and PSO algorithms in reverse engineering using metaheuristic algorithms is confirmed.

## 5. Conclusions

This study proposes a new method for the RE of helical gears. This issue has been less addressed in previous studies in the literature. The main objective of the study was to find accurate gear design parameters of a helical gear, such as the normal module $m_{n}$, normal pressure angle $\alpha_{n}$, addendum modification $x$, and helix angle $\beta$, using performed span and over-balls measurements. The relation between the input design parameters and relevant output measurable characteristics is obtainable with complex mathematical equations. Since it is impossible to mathematically solve the mentioned inverse equations, metaheuristic methods should be used. First, an objective function based on measurements and equations was defined. Then the defined problem was optimized by four different algorithms that have widespread use in engineering problems: GWO, WOA, PSO, and GA. The convergence speed and stability of the proposed methods were compared via statistical measures, and various indicators of the algorithms such as the convergence speed and number of function evaluations were examined. Finally, the following results were obtained:

- Based on the proposed methodology, accurate input parameters were reached. The validation of the obtained results was evaluated by the given equations.
- The performance of the algorithms was assessed in terms of the ability to reach the best solution, convergence speed, and stability. It was found that swarm-based optimization methods such as GWO and PSO are superior to the other considered algorithms such as GA.
For future work, applying the proposed methodology in the RE of other types of gears such as bevel gears and worm gears is suggested. Examining the other natureinspired algorithms and recently published and advanced optimization methods in the RE of mechanical parts, especially gear components, is another field of research that could be considered extensively. Another vast area of research is applying the various discussed tuners in metaheuristics and their effect on the convergence speed of algorithms.

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# Designing a Secure Mechanism for Image Transferring System Based on Uncertain Fractional Order Chaotic Systems and NLFPID Sliding Mode Controller 

Mohammad Rasouli ${ }^{1, *(\mathbb{D}}$, Assef Zare ${ }^{2,3, *(\mathbb{D}}$, Hassan Yaghoubi ${ }^{4}$ (D) and Roohallah Alizadehsani ${ }^{5}$ (D)<br>1 College of Skills and Entrepreneurship, Mashhad Branch, Islamic Azad University, Mashhad 9187147578, Iran<br>2 Faculty of Electrical Engineering, Gonabad Branch, Islamic Azad University, Gonabad 6518115743, Iran<br>3 Research Center of Intelligent Technologies in Electrical Industry, Islamic Azad University, Gonabad 6518115743, Iran<br>4 Department of Electrical and Electronic Engineering, Gonabad Branch, Islamic Azad University, Gonabad 6518115743, Iran; yaghoubihassan@yahoo.com<br>5 Institute for Intelligent Systems Research and Innovation (IISRI), Deakin University, Waurn Ponds, Geelong, VIC 3216, Australia; r.alizadehsani@deakin.edu.au<br>* Correspondence: mohamad1rasouli@gmail.com (M.R.); assefzare@gmail.com (A.Z.)

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#### Abstract

A control method for the robust synchronization of a class of chaotic systems with unknown time delay, unknown uncertainty, and unknown disturbance is presented. The robust controller was designed using a nonlinear fractional order PID sliding surface. The Lyapunov method was used to determine the update laws, prove the stability of the proposed mechanism, and guarantee the convergence of the synchronization errors to zero. The simulation was performed using MATLAB software to evaluate the performance of the proposed mechanism, and the results showed that it was efficient. Finally, the proposed method was combined with a secure communication application to encrypt images, and the results obtained were favorable regarding the standard criteria of correlation, NPCR, PSNR, and information entropy.


Keywords: chaotic synchronization; fractional order sliding mode control; adaptive control; secure communication

MSC: 93D09; 93B51

## 1. Introduction

Chaos is a nonlinear phenomenon that appears to be random but actually follows a pattern. It was discovered about a half-century ago by Lorenz [1]. Scientists began to pay more attention to the phenomenon of chaos after that. Some systems, including the Liu system [2], the IU system [3], and the Chen system [4], have been proposed on the basis of Lorenz's ideas. About 300 years ago, fractional calculations were introduced, and more complete definitions and theorems have been introduced since then [5]. Physical systems can be represented as integer or fractional equations in this context. It is evident that modeling using fractional order systems can have more accuracy than modeling with integer order systems. Recently, the description of systems using fractional calculus has been developed in various sciences, including chemical reaction systems [6,7], biological systems [8,9], power converters [10], electrochemical processes [11], robotics [12], and others. The problem of synchronizing two chaotic systems has piqued the interest of scientists working in the field of secure communication over the past two decades. In fact, the synchronization of two chaotic systems can be described as a situation in which two or more chaotic systems coordinate their responses by the controller. As a result, two subsystems, the main or driving system and the slave or response system, constitute
a coupled system. The master system's response is unconstrained and drives the slave system.

To tackle the problem, Petras et al. presented a fractional sliding surface [13]. Zare Hallaji et al. [14,15] presented research on the synchronization of positive and fractional chaotic systems with system uncertainty. They evaluated the conditions of the described problem from several perspectives, including unknown uncertainties in the system characteristics, in their research. In [16], the chaotic system was synchronized using a nonlinear observer and the benefit of adaptive control in order to determine the system's uncertainties. In this design, a sliding surface equivalent to one of the system states was provided, and its stability was demonstrated using a Lyapunov function. The authors of [17] proposed an adaptive terminal sliding mode controller (ATSMC). First, a fractional order sliding surface for the master and slave system was introduced in this article. The stability of the suggested controller was then examined, as was the ongoing convergence of the error in the synchronization problem.

A sliding surface based on the nonlinear fractional order PID was developed in this study for the synchronization of two systems with uncertainty and unknown disturbances with unknown and time-varying time delay. The following benefits might be highlighted in this research, which was conducted to synchronize two systems:

- The use of the nonlinear fractional PID (NLFOPID) sliding surface instead of typical sliding surfaces.
- The presence of unknown time delays
- The presence of uncertainty and disturbance with unknown boundaries. Then, using the suitable Lyapunov function and update laws, a control signal was extracted that could be used to overcome the chattering problem by properly adjusting the controller parameters. This is a critical issue for the suggested controller's implementation. In $[10,18]$, a controller for the synchronization of chaotic systems in finite time was constructed utilizing a sliding surface, and the synchronization of the integer order chaotic system was investigated in [19].
The preliminary calculations of deficit accounts are reported in Section 2 of this article. Section 3 presents the equations characterizing the system as well as the set limitations for uncertainty. Section 4 introduces the sliding surface based on the proportional-integralnonlinear fractional derivative, as well as the controller architecture. Section 5 investigates the adaptive controller's stability analysis and update laws. Section 6 presents the simulation results and visualization of the synchronized system. Section 7 discusses chaotic masking for image encryption. Finally, in the last section, conclusions and recommendations are offered.


## 2. Preliminary Definitions of Fractional Order Differentiation

Definition 1. The fractional order integration and differentiation are defined as follows [20]:

$$
D_{t}^{q}=\left\{\begin{array}{l}
\frac{d^{q}}{d t^{q}} \quad q>0  \tag{1}\\
1 q=0 \\
\int_{q}^{t}(d \tau)^{-q} \quad q<0
\end{array}\right.
$$

in which $q$ is a real number.
Definition 2. The Riemann-Liouville fractional integral of order $q$ of the function $f(t)$ is defined as follows [21]:

$$
\begin{equation*}
t_{0} I_{t}^{q} f(t)=\frac{1}{\Gamma(q)} \int_{t_{0}}^{t} \frac{f(\tau)}{(t-\tau)^{1-q}} d \tau \tag{2}
\end{equation*}
$$

in which $t_{0}$ is the initial time and $\Gamma(q)$ is the Gamma function defined as follows:

Definition 3. Suppose $n-1<q \leq n, n \in N$. The fractional Riemann-Liouville differentiation of order $q$ is defined for the function $f(t)$ below [21]:

$$
\begin{equation*}
t_{0} D_{t}^{q} f(t)=\frac{d^{q} f(t)}{d t^{q}}=\frac{1}{\Gamma(n-q)} \frac{d^{n}}{d t^{n}} \int_{t_{0}}^{t} \frac{f(\tau)}{(t-\tau)^{q-n+1}} d \tau \tag{3}
\end{equation*}
$$

Note 1: In Equation (4), the Riemann-Liouville fractional order integral is first calculated, and then differentiation is performed; thus, the derivative of a constant number in this formulation is not equal to zero.

Definition 4. In the continuous function $f(t)$, the Caputo fractional order derivative of order $q$ is defined as follows [21]:

$$
\begin{gather*}
t_{0} D_{t}^{q} f(t)=\left\{\begin{array}{cc}
\frac{1}{\Gamma(m-q)} \int_{t_{0}}^{t} \frac{f^{(m)}(\tau)}{\left(t-\tau \tau^{q-m+1}\right.} d \tau & m-1<q<m \\
\frac{d^{m} f(t)}{d t^{m}} & q=m
\end{array}\right.  \tag{4}\\
\Gamma(q)=\int_{0}^{\infty} e^{-t} t^{q-1} d t \tag{5}
\end{gather*}
$$

Such that $m$ is the first integer number after $q$.
Lemma 1. If $f(t)$ is a constant function and $q>0$, the Caputo derivative in Equation (5) for $f(t)$ would be as follows:

$$
\begin{equation*}
D^{q} f(t)=0 \tag{6}
\end{equation*}
$$

The authors of [22] presented the stability analysis of fractional order systems using the direct Lyapunov method, as well as the determination of the necessary and sufficient conditions guaranteeing stability using the Mittag-Leffler concept, and the authors of [23] reviewed the stability analysis of nonlinear systems using convex Lyapunov functions.

Lemma 2 [23]. Suppose that $h(t) \in R$ is a continuous and differentiable function. Then, for $t \geq t_{0}$, Equation (7) is satisfied.

$$
\begin{equation*}
D^{q} h^{2}(t) \leq 2 h(t) \cdot D^{q} h(t) \tag{7}
\end{equation*}
$$

Lemma 3 [23]. Suppose that $h(t) \in R^{n}$ is a continuous and differentiable function. Then, for $t \geq t_{0}$, we have:

$$
\begin{equation*}
D^{q} h^{T}(t) \cdot h(t) \leq 2 h^{T}(t) \cdot D^{q} h(t) \tag{8}
\end{equation*}
$$

Theorem 1 [22]. Assume that the origin $(x=0)$ is the equilibrium point of the fractional order system (5) and that its definition domain covers the origin. Furthermore, $v(x(t), t)$ is a continuous and differentiable Lipschitz function, implying the following:

$$
\begin{align*}
& D^{q} x(t)=f(x, t) \\
& a_{1}\|x\|^{a} \leq v(x(t), t) \leq a_{2}\|x\|^{a b}  \tag{9}\\
& D^{q_{v}}(x(t), t) \leq-a_{3}\|x\|^{a b}
\end{align*}
$$

in which $0<q<1$ and $a, a_{1}, a_{2}, a_{3}, b$ are positive arbitrary constants. Then, the origin is stable in the Mittag-Leffler sense.

Definition 5. The continuous function $p:[0, \infty) \rightarrow[0, \infty)$ belongs to class $k$ if its derivative is positive and $p(0)=0$.

Theorem 2 [22]. Assume $x=0$ is the equilibrium point of the fractional order system (5), the Lipschitz condition for $f(x, t)$ is satisfied, and $q \in(0,1)$. If Equations (8) and (9) are satisfied for the Lyapunov function $v(x(t), t)$ and functions $\delta_{i}$ of class $K$ :

$$
\begin{align*}
& \delta_{1}(\|x\|) \leq v(x(t), t) \leq \delta_{2}(\|x\|)  \tag{10}\\
& D^{q} v(x(t), t) \leq-\delta_{3}(\|x\|)
\end{align*}
$$

Then, system (5) is asymptotically stable in the Mittag-Leffler sense.
Theorem 3 [24]. For the fractional order system (5) and the Lyapunov function $v(x)$, we have:

$$
\begin{equation*}
D^{q} v(x) \leq\left(\frac{\partial v}{\partial x}\right)^{T} \cdot D^{q} x=\left(\frac{\partial v}{\partial x}\right)^{T} \cdot f(x, t) \tag{11}
\end{equation*}
$$

Definition 6 [25]. A continuous piecewise function $f(x, t)$ has the Lipschitz condition if:

$$
\begin{equation*}
\|f(x, t)-f(z, t)\| \leq \gamma_{f}\|x-z\|, \forall x, z \in R^{n} \tag{12}
\end{equation*}
$$

## 3. System Descriptor Equations

The equations characterizing a class of master-slave chaotic systems with uncertainty and indeterminate time delay in the presence of an unknown disturbance are introduced in this section. Following standardization, the master system dynamics in canonical form are as follows:

$$
\left\{\begin{array}{c}
D^{q} x_{i}=x_{i+1} 1 \leq i \leq n-1  \tag{13}\\
D^{q} x_{n}=\sigma_{0}^{T} x+f\left(x\left(t-\tau_{1}\right), t\right)+\Delta f(x(t), t)+d_{1}(t)
\end{array}\right.
$$

The slave system equations are as follows:

$$
\left\{\begin{array}{c}
D^{q} y_{i}=y_{i+1} 1 \leq i \leq n-1  \tag{14}\\
D^{q} y_{n}=\sigma_{0}^{T} y+g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(y(t), t)+d_{2}(t)+u(t)
\end{array}\right.
$$

The differential equations are written in the forms of well-known chaotic systems, such as the Van der Pol Oscillator, Duffing's Oscillator, the Genesio-Tesi System, Arneodo's System, and so on [26], where $x(t), y(t) \in R^{n}$ denote the dynamic states of the master and slave systems, $\sigma_{0}^{T}$ denotes the constant coefficients in the system's linear states, and $f\left(x\left(t-\tau_{1}\right), t\right), g\left(y\left(t-\tau_{2}\right), t\right) \in R$ are nonlinear functions with an unknown delay with $\tau_{1}, \tau_{2}$ delays, and $\Delta f(x(t), t), \Delta g(x(t), t)$ represent bounded uncertainty in the master and slave systems. Furthermore, $d_{1}(t), d_{2}(t)$ indicate the external distortions applied to the master and slave systems, respectively, while $u(t)$ is the control law applied to the slave system.

Definition 7. If the following conditions are satisfied for the systems described in Equations (13) and (14) for all the conditions governing the system, including all initial conditions, uncertainties, unknown time delay, and external disturbance, the system has robust synchronization:

$$
\begin{equation*}
\lim _{t \rightarrow \infty}\left|y_{i}(t)-x_{i}(t)\right|=\lim _{t \rightarrow \infty}\left|e_{i}(t)\right|=0, i=1, \ldots, n \tag{15}
\end{equation*}
$$

As a result, $e_{i}(t)$ introduces the synchronization error of the master and slave systems.
As a result, the following are the dynamic equations describing the synchronization error for the uncertain chaotic master and slave systems with unknown time delay described in (13) and (14):

$$
\left\{\begin{array}{c}
D^{q} e_{i}=e_{i+1} 1 \leq i \leq n-1  \tag{16}\\
D^{q} e_{n}=\sigma_{0}^{T}(\mathrm{y}-\mathrm{x})+g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t) \\
-\left(f\left(x\left(t-\tau_{1}\right), t\right)+\Delta f(x(t), t)+d_{1}(t)\right)+u(t) .
\end{array}\right.
$$

Assumption 1. The uncertain external disturbances $\boldsymbol{d}_{\mathbf{1}}(\boldsymbol{t}) \cdot \boldsymbol{d}_{\mathbf{2}}(\boldsymbol{t})$ and the uncertain bounded nonlinear uncertainties $\Delta f(x(t) \cdot t)$ and $\Delta g(x(t) \cdot \boldsymbol{t})$ in the master and slave systems (13) and (14) meet the following conditions:

$$
\begin{align*}
& \|\Delta f(x(t), t)\| \leq \beta_{1} \omega_{1}(x) \\
& \|\Delta \boldsymbol{g}(\boldsymbol{y}(t), t)\| \leq \beta_{2} \omega_{2}(y) \\
& \left\|\boldsymbol{d}_{\mathbf{1}}(\boldsymbol{t})\right\| \leq \rho_{\mathbf{1}}  \tag{17}\\
& \left\|\boldsymbol{d}_{\mathbf{2}}(\boldsymbol{t})\right\| \leq \boldsymbol{\rho}_{\mathbf{2}} \\
& \underline{\boldsymbol{\tau}}_{\mathbf{i}}<\boldsymbol{\tau}_{\mathbf{i}}<\boldsymbol{\tau}_{\mathbf{i}}
\end{align*}
$$

Such that $\|$.$\| denotes the l_{1}$ norm, $\beta_{2}, \beta_{1}, \rho_{2}, \rho_{1}$ are unknown positive real numbers, and $\omega_{2}(\cdot), \omega_{1}(\cdot)$ are positive and known functions. Also, $\rho_{i}<\bar{\rho}_{i}, \beta_{i}<\bar{\beta}_{i}$ where $\bar{\rho}_{i}, \bar{\beta}_{i}, \tau_{i}$, and $\underline{\tau}_{i}$ are known values.
Assumption 2. The nonlinear functions $f\left(x\left(t-\tau_{1}\right), t\right), g\left(y\left(t-\tau_{2}\right), t\right) \in R$ satisfy the Lipschitz conditions for any $x(t), y(t) \in R$ :

$$
\begin{align*}
& \left|f\left(x\left(t-\tau_{1}\right)\right)-f\left(x\left(t-\hat{\tau}_{1}\right)\right)\right| \leq l_{1}\left|\tau_{1}-\hat{\tau}_{1}\right|=l_{1}\left|\tilde{\tau}_{1}\right|  \tag{18}\\
& \left|g\left(y\left(t-\tau_{2}\right)\right)-g\left(y\left(t-\hat{\tau}_{2}\right)\right)\right| \leq l_{2}\left|\tau_{2}-\hat{\tau}_{2}\right|=l_{2}\left|\tilde{\tau}_{2}\right|
\end{align*}
$$

Table 1 presents the system parameters and the proposed mechanism:
Table 1. Symbols and concepts.

| Symbol | Concept | Symbol | Concept |
| :---: | :--- | :---: | :--- |
| $\rho_{i}$ | Disturbance bound | $\hat{\rho}_{i}$ | Disturbance bound estimate |
| $\beta_{i}$ | Uncertainty bound | $\hat{\beta}_{i}$ | Uncertainty bound estimate |
| $l_{i}$ | Lipschitz constant | $\hat{\tau}_{i}$ | Time delay bound estimate |
| $\tau_{i}$ | Time delay | $\tilde{\rho}_{i}$ | Disturbance bound estimate error |
| $\bar{\rho}_{i}$ | Disturbance upper bound | $\tilde{\beta}_{i}$ | Uncertainty bound estimate error |
| $\bar{\beta}_{i}$ | Uncertainty upper bound | $\tilde{\tau}_{i}$ | Time delay estimate error |
| $\tau_{i}$ | Time delay upper bound | $b$ | Positive constant number |
| $\tau_{i}$ | Time delay lower bound | $\bar{\epsilon}$ | Small positive constant number |

In this study, all states of the system were directed to and kept on the sliding surface by designing a robust adaptive controller and introducing an integral proportional sliding surface and a fractional order nonlinear derivative. Furthermore, the system's uncertainties and unknown parameters should be estimated and updated. Then, in the robust synchronization of chaotic systems (13) and (14) in the presence of external distortions, bounded nonlinear uncertainties, and uncertain time delays, the dynamics of the slave system state must match the behavior of the master system dynamics, and the estimation error of the unknown parameters in both chaotic systems approach zero in any circumstance, ensuring the system's robust stability.

## 4. The Sliding Mode Control Approach Based on Fractional Order Nonlinear PID Controllers

A proportional integral sliding surface and a nonlinear fractional order derivative are presented in this section in order to synchronize chaotic systems (13) and (14) with unknown uncertainty and unknown time delay. The fractional order sliding surface is as follows, according to the nonlinear fractional order PID controller structure presented in [26], which enhances tracking:

$$
\begin{equation*}
s(t)=h(e) \cdot\left[k_{p} e_{n}(t)+T_{I} D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}+T_{d} D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right] \tag{19}
\end{equation*}
$$

Such that $h(e)$ is a nonlinear function, defined as follows:

$$
\begin{equation*}
h(e)=k_{0}+\left(1-k_{0}\right)\|E(t)\|, k_{0} \in(0,1) \tag{20}
\end{equation*}
$$

where $\|E(t)\|=\sum_{i=1}^{n}\left|e_{i}\right|$. Coefficients $T_{I}$ and $T_{d}$ are time constants of integral and derivative sentences. The parameters $k_{1 i}$ and $k_{2 i}$ are positive constant values of the sliding surface such that they satisfy the stability of the desired system. If the system is in sliding mode, the following conditions must be met:

$$
\begin{equation*}
s(t)=0, D^{q} s(t)=0 \tag{21}
\end{equation*}
$$

The fractional order derivative of the sliding surface in Equation (21) is as follows:

$$
\begin{align*}
D^{q} s(t)= & \left(k_{0} k_{p} D^{q} e_{n}(t)+k_{0} T_{i} D^{q-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{q+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)+\left(1-k_{0}\right) k_{p} D^{q}\left(\|E(t)\| e_{n}(t)\right)\right.  \tag{22}\\
& \left.+\left(1-k_{0}\right) T_{I} D^{q}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right)+\left(1-k_{0}\right) T_{d} D^{q}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right)=0
\end{align*}
$$

Now, $D^{q} e_{n}$ is substituted into Equation (21) using Equation (16):

$$
\begin{align*}
D^{q} s(t)= & \left(k _ { 0 } k _ { p } \left(g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t)-\left(f\left(x\left(t-\tau_{1}\right), t\right)+\Delta f(x(t), t)+d_{1}(t)\right)+\sigma_{0}^{T} \cdot E(t)\right.\right. \\
& +u(t))+k_{0} T_{i} D^{1-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{1+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)+\left(1-k_{0}\right) k_{p} D^{q}\left(\|E(t)\| e_{n}(t)\right)  \tag{23}\\
& \left.+\left(1-k_{0}\right) T_{I} D^{q}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right)+\left(1-k_{0}\right) T_{d} D^{q}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right)=0
\end{align*}
$$

In this case, the control signal is determined as follows:

$$
\begin{align*}
u(t)= & \frac{-1}{k_{0} k_{p}}\left(k_{0} T_{i} D^{q-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{q+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)+\left(1-k_{0}\right) k_{p} D^{q}\left(\|E(t)\| e_{n}(t)\right)\right. \\
& \left.+\left(1-k_{0}\right) T_{I} D^{q}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right)+\left(1-k_{0}\right) T_{d} D^{q}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right)  \tag{24}\\
& +f\left(x\left(t-\hat{\tau}_{1}\right) \cdot t\right)-g\left(y\left(t-\hat{\tau}_{2}\right) \cdot t\right)-\sigma_{0}^{T} \cdot E(t)-b s+\bar{u}(t)
\end{align*}
$$

In Equation (25), the term $\bar{u}(t)$ comprises the terms coming from the estimation of the system's bounds of uncertainties and disturbances, which are defined using the adaptive controller, as follows:

$$
\begin{align*}
& \left.\bar{u}(t)=-\operatorname{sgn}(s)\left[\hat{\beta}_{2} \omega_{2}(\mathbf{y})+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right)\right]+u_{00}(t) \\
& u_{00}(t)=\frac{-b}{k_{0} k_{p} s} \sum_{i=1}^{2}\left[\left(\left|\hat{\rho}_{i}\right|+\bar{\rho}_{i}\right)^{2}+\left(\left|\hat{\tau}_{i}\right|+\bar{\tau}_{i}\right)^{2}+\left(\left|\hat{\beta}_{i}\right|+\bar{\beta}_{i}\right)^{2}\right] \tag{25}
\end{align*}
$$

## 5. Stability Analysis and Determining the Update Laws

The construction of the robust adaptive controller is described in this part, employing the sliding surface based on nonlinear fractional order PID in such a way that the suggested control strategy guarantees the stability of the synchronization of chaotic systems.

Theorem 4. The synchronization of systems (13) and (14) in the presence of disturbances $d_{1}$ and $d_{2}$ and unknown uncertainties $\Delta f$ and $\Delta g$ with unknown time delays $\tau_{1}$ and $\tau_{2}$ and the definition of the controller $u(t)$ is guaranteed as follows:

$$
\begin{align*}
u(t)= & -g\left(y\left(t-\hat{\tau}_{1}\right)\right)+f\left(x\left(t-\hat{\tau}_{2}\right)\right) \\
& -\frac{1}{k_{0} k_{p}}\left(k_{0} T_{l} D^{q-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{q+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)+\left(1-k_{0}\right) k_{p} D^{q}\left(\|E(t)\| e_{n}(t)\right)\right.  \tag{26}\\
& \left.+\left(1-k_{0}\right) T_{I} D^{q}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right)+\left(1-k_{0}\right) T_{d} D^{q}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)\right)-\sigma_{0}^{T} \\
& \cdot E(t)-b s-\operatorname{sgn}(s)\left(\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right)+u_{00}(t)
\end{align*}
$$

Such that the update laws are as follows:

$$
\begin{align*}
& D^{q} \hat{\tau}_{i}=-D^{q} \tilde{\tau}_{i}=l_{i}|s| \operatorname{sgn}\left(\tilde{\tau}_{i}\right), \hat{\tau}_{i}(0)=\bar{\tau}_{i} \\
& D^{q} \hat{\rho}_{i}=-D^{q} \tilde{\rho}_{i}=k_{0} k_{p}|s|  \tag{27}\\
& D^{q} \hat{\beta}_{1}=-D^{q} \tilde{\beta}_{1}=-k_{0} k_{p}|s| \omega_{2}(y) \\
& D^{q} \hat{\beta}_{2}=-D^{q} \tilde{\beta}_{2}=-k_{0} k_{p}|s| \omega_{1}(x)
\end{align*}
$$

Thus, the convergence of the chaotic systems' synchronization error to zero is ensured.
Proof. Consider the following Lyapunov function:

$$
\begin{equation*}
v(t)=\frac{1}{2}\left[s^{2}(t)+\stackrel{\sim}{\beta}_{1}^{2}+\sim_{\beta}^{2}+l_{1} \tilde{\tau}_{1}^{2}+l_{2} \tilde{\tau}_{2}^{2}+\tilde{\rho}_{1}^{2}+\tilde{\rho}_{2}^{2}\right] \tag{28}
\end{equation*}
$$

in which the parameters' estimation error is defined as follows:

$$
\begin{equation*}
\tilde{\tau}_{i}=\tau_{i}-\hat{\tau}_{i}, \tilde{\rho}_{i}=\rho_{i}-\hat{\rho}_{i}, \tilde{\beta}_{i}=\beta_{i}-\hat{\beta}_{i} \tag{29}
\end{equation*}
$$

Considering Equation (28), the derivative of the Lyapunov function is as follows:

$$
\begin{equation*}
D^{q} v(t)=\frac{1}{2} D^{q}\left(s^{2}+\tilde{\beta}_{1}^{2}+\tilde{\beta}_{2}^{2}+l_{1} \tilde{\tau}_{1}^{2}+l_{2} \tilde{\tau}_{2}^{2}+\tilde{\rho}_{1}^{2}+\tilde{\rho}_{2}^{2}\right) \leq s \cdot D^{q} s+\sum_{i=1}^{2}\left(\tilde{\beta}_{i} D^{q} \tilde{\beta}_{i}+l_{i} \tilde{\tau}_{i} D^{q} \tilde{\tau}_{i}+\tilde{\rho}_{i} D^{q} \tilde{\rho}_{i}\right) \tag{30}
\end{equation*}
$$

By applying Equation (23) in Equation (30), Equation (31) is determined:

$$
\begin{array}{rl}
D^{q} v(t) \leq s & s \cdot\left[k _ { 0 } k _ { p } \left(g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t)-\left(f\left(x\left(t-\tau_{1}\right), t\right)+\Delta f(x(t), t)+d_{1}(t)\right)+\sigma_{0}^{T} \cdot E(t)\right.\right. \\
& +u(t))+k_{0} T_{i} D^{q-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)+k_{0} T_{d} D^{q+\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)+\left(1-k_{0}\right) k_{p} D^{q}\left(\|E(t)\| e_{n}(t)\right) \\
& \left.+\left(1-k_{0}\right) T_{I} D^{q}\left(\|E(t)\| D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}(t)\right)+\left(1-k_{0}\right) T_{d} D^{q}\left(\|E(t)\| D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)\right)+u_{00}(t)\right]  \tag{31}\\
& +\sum_{i=1}^{2}\left(\tilde{\beta}_{i} D^{q} \tilde{\beta}_{i}+l_{i} \tilde{\tau}_{i} D^{q} \tilde{\tau}_{i}+\tilde{\rho}_{i} D^{q} \tilde{\rho}_{i}\right)
\end{array}
$$

In this case, the Lyapunov function derivate is as follows:

$$
\begin{align*}
& D^{q} v(t) \leq s \cdot\left[k _ { 0 } k _ { p } \left(g\left(y\left(t-\tau_{2}\right), t\right)-g\left(y\left(t-\hat{\tau}_{2}\right), t\right)+\Delta g(x(t), t)+d_{2}(t)+f\left(x\left(t-\hat{\tau}_{1}\right), t\right)-f\left(x\left(t-\tau_{1}\right), t\right)\right.\right. \\
& \left.\left.-\Delta f(x(t), t)-d_{1}(t)-b s-\operatorname{sgn}(s)\left[\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right]\right)\right]+s k_{0} k_{p} u_{00}(t)  \tag{32}\\
& +\sum_{i=1}^{2}\left(\tilde{\beta}_{i} D^{q} \tilde{\beta}_{i}+l_{i} \tilde{\tau}_{i} D^{q} \tilde{\tau}_{i}+\tilde{\rho}_{i} D^{q} \tilde{\rho}_{i}\right)
\end{align*}
$$

Thus, we have:

$$
\begin{align*}
& D^{q} v(t) \leq \quad|s| \cdot\left[k _ { 0 } k _ { p } \left(\left|g\left(y\left(t-\tau_{2}\right), t\right)-g\left(y\left(t-\hat{\tau}_{2}\right), t\right)\right|+|\Delta g(x(t), t)|+\left|f\left(x\left(t-\hat{\tau}_{1}\right), t\right)-f\left(x\left(t-\tau_{1}\right), t\right)\right|\right.\right. \\
&\left.\left.-|\Delta f(x(t), t)|+\left|d_{2}(t)-d_{1}(t)\right|\right)\right]-k_{0} k_{p} b s^{2} \\
&+k_{0} k_{p} s\left(-\operatorname{sgn}(s)\left[\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right]\right)+s k_{0} k_{p} u_{00}(t)  \tag{33}\\
&+\sum_{i=1}^{2}\left(\tilde{\beta}_{i} D^{q} \tilde{\beta}_{i}+l_{i} \tilde{\tau}_{i} D^{q} \tilde{\tau}_{i}+\tilde{\rho}_{i} D^{q} \tilde{\rho}_{i}\right)
\end{align*}
$$

On the basis of assumptions 1-2 and 2-2 presented in Equations (17) and (18) in Section 3 of the article, Equation (33) is rewritten as follows:

$$
\begin{align*}
& D^{q} v(t) \leq \quad|s| \cdot\left[k_{0} k_{p}\left(l_{2}\left|\tau_{2}-\hat{\tau}_{2}\right|+\beta_{2} \omega_{2}(y)+l_{1}\left|\tau_{1}-\hat{\tau}_{1}\right|+\beta_{1} \omega_{1}(x)+\rho_{1}+\rho_{2}\right)\right]-k_{0} k_{p} b s^{2} \\
&-k_{0} k_{p} \operatorname{sgn}(s) s\left[\hat{\beta}_{2} \omega_{2}(y)+\hat{\beta}_{1} \omega_{1}(x)+\hat{\rho}_{2}+\hat{\rho}_{1}\right]+s k_{0} k_{p} u_{00}(t)  \tag{34}\\
&+\sum_{i=1}^{2}\left(\tilde{\beta}_{i} D^{q} \tilde{\beta}_{i}+l_{i} \tilde{\tau}_{i} D^{q} \tilde{\tau}_{i}+\tilde{\rho}_{i} D^{q} \tilde{\rho}_{i}\right)
\end{align*}
$$

The derivative of the Lyapunov function is as follows:

$$
\begin{align*}
D^{q} v(t) \leq & |s|\left[k_{0} k_{p}\left(l_{1}\left|\tilde{\tau}_{1}\right|+\tilde{\beta}_{2} \omega_{2}(y)+l_{2}\left|\tilde{\tau}_{2}\right|+\tilde{\beta}_{1} \omega_{1}(x)+\tilde{\rho}_{2}+\tilde{\rho}_{1}\right)\right]-b s^{2}+s k_{0} k_{p} u_{00}(t)  \tag{35}\\
& +\sum_{i=1}^{2}\left(\tilde{\beta}_{i} D^{q} \tilde{\beta}_{i}+l_{i} \tilde{\tau}_{i} D^{q^{q}} \tilde{\tau}_{i}+\tilde{\rho}_{i} D^{q} \tilde{\rho}_{i}\right)
\end{align*}
$$

Now, by substituting the update laws (27) into (35), the derivative of the Lyapunov function is simplified as follows:

$$
\begin{equation*}
\Rightarrow D^{q} v(t) \leq-b s^{2}+s k_{0} k_{p} u_{00}(t) \tag{36}
\end{equation*}
$$

In the following, by substituting $u_{00}(t)$ from (25) into (36), Equation (37) is obtained:

$$
\begin{equation*}
\Rightarrow D^{q} v(t) \leq-b s^{2}-s k_{0} k_{p} \frac{b}{k_{0} k_{p} s} \sum_{i=1}^{2}\left[\left(\left|\hat{\rho}_{i}\right|+\bar{\rho}_{i}\right)^{2}+\left(\left|\hat{\tau}_{i}\right|+\bar{\tau}_{i}\right)^{2}+\left(\left|\hat{\beta}_{i}\right|+\bar{\beta}_{i}\right)^{2}\right] \tag{37}
\end{equation*}
$$

On the other hand:

$$
\begin{align*}
& \left|\tilde{\tau}_{i}\right|=\left|\tau_{i}-\hat{\tau}_{i}\right| \leq\left|\tau_{i}\right|+\left|\hat{\tau}_{i}\right| \leq\left|\hat{\tau}_{i}\right|+\bar{\tau}_{i} \Rightarrow-\left(\left|\hat{\tau}_{i}\right|+\bar{\tau}_{i}\right)^{2} \leq-\left|\tilde{\tau}_{i}\right|^{2}  \tag{38}\\
& \left|\tilde{\beta}_{i}\right|=\left|\beta_{i}-\hat{\beta}_{i}\right| \leq\left|\beta_{i}\right|+\left|\hat{\beta}_{i}\right| \leq\left|\hat{\beta}_{i}\right|+\bar{\beta}_{i} \Rightarrow-\left(\left|\hat{\beta}_{i}\right|+\bar{\beta}_{i}\right)^{2} \leq-\left|\tilde{\beta}_{i}\right|^{2} \\
& \left|\tilde{\rho}_{i}\right|=\left|\rho_{i}-\hat{\rho}_{i}\right| \leq\left|\rho_{i}\right|+\left|\hat{\rho}_{i}\right| \leq\left|\hat{\rho}_{i}\right|+\bar{\rho}_{i} \Rightarrow-\left(\left|\hat{\rho}_{i}\right|+\bar{\rho}_{i}\right)^{2} \leq-\left|\tilde{\rho}_{i}\right|^{2}
\end{align*}
$$

By substituting Equation (38) into Equation (35), the derivative of the Lyapunov function is simplified to Equation (39).

$$
\begin{equation*}
\Rightarrow D^{q} v(t) \leq-b\left(s^{2}+\sum_{i=1}^{2}\left[\tilde{\beta}_{i}^{2}+\tilde{\tau}_{i}^{2}+\tilde{\rho}_{i}^{2}\right]\right) \leq-2 b v \tag{39}
\end{equation*}
$$

The convergence of $v(t)$ to zero is guaranteed by Theorems (1) and (2). As a result, the sliding surface $s$ and the estimation errors approach zero. In the following, it is proven that the synchronization errors approach zero. For this purpose, first, $\alpha_{i} \triangleq T_{I} k_{1 i}$ and $\beta_{i} \triangleq T_{d} k_{2 i}$ are defined. Then, by applying Equations (19)-(21), expression (40) is obtained:

$$
\begin{equation*}
\Rightarrow k_{p} e_{n}(t)+T_{I} D^{-\lambda} \sum_{i=1}^{n} k_{1 i} e_{i}+T_{d} D^{\delta} \sum_{i=1}^{n} k_{2 i} e_{i}(t)=0 \tag{40}
\end{equation*}
$$

Thus, the fractional order derivative of is obtained from both sides of Equation (37)

$$
\begin{equation*}
k_{p} D^{\lambda} e_{n}(t)+\sum_{i=1}^{n} \beta_{i} e_{i}(t)+\sum_{i=1}^{n} \alpha_{i} D^{\lambda+\delta} e_{i}(t)=0,0<\lambda+\delta \leq 1 \tag{41}
\end{equation*}
$$

The dynamics of the system error are defined as follows:

$$
\left\{\begin{array} { c } 
{ D ^ { q } e _ { 1 } = e _ { 2 } }  \tag{42}\\
{ D ^ { q } e _ { 2 } = e _ { 3 } } \\
{ \vdots } \\
{ D ^ { q } e _ { n - 1 } = e _ { n } }
\end{array} \Rightarrow \left\{\begin{array}{c}
s^{q} E_{1}=E_{2}+k_{2}^{\prime}(s) \\
s^{q} E_{2}=E_{3}+k_{3}^{\prime}(s) \\
\vdots \\
s^{q} E_{n-1}=E_{n}+k_{n}^{\prime}(s)
\end{array} \Rightarrow E_{i}=s^{(i-1) q} E_{1}(s)+k_{i}^{\prime}(s)\right.\right.
$$

in which $E_{i}(s)=\mathcal{L}\left(e_{i}\right)$, and $k_{i}^{\prime}(s)$ is the effect of the initial condition of the Laplace transform. By calculating the Laplace transform using Equation (40), Equation (42) is obtained:

$$
\begin{equation*}
k_{p} s^{\lambda} E_{n}(s)+\sum_{i=1}^{n}\left(\alpha_{i} s^{q+\lambda} E_{i}+\beta_{i} E_{i}\right)=k_{0}(s) \tag{43}
\end{equation*}
$$

where $k_{0}(s)$ is the general effect of the initial conditions. By substituting Equation (42) into Equation (43), Equation (44) is obtained:

$$
\begin{equation*}
\left[k_{p} s^{\lambda} s^{(n-1) q}+\sum_{i=1}^{n}\left(\alpha_{i} s^{\delta+\lambda_{s}} s^{(i-1) q}+\beta_{i} s^{(i-1) q}\right)\right] E_{1}(s)=k_{0}(s) \tag{44}
\end{equation*}
$$

Therefore, the system's characteristic equation is as follows:

$$
\begin{equation*}
k_{p} s^{(n-1) q+\lambda}+\sum_{i=1}^{n}\left(\alpha_{i} s^{\delta+\lambda+(i-1) q}+\beta_{i} s^{(i-1) q}\right)=0 \tag{45}
\end{equation*}
$$

If the coefficients $\alpha_{i}, \beta_{i}$, and $k_{p}$ on the sliding surface are chosen in such a way that the roots of the above equation have a negative real part, then all $e_{i} \mathrm{~s}$ approach zero.

Therefore, a sufficient condition for the synchronization errors to converge to zero is that the characteristic Equation (45) is stable.

In Equation (25), if the sliding surface approaches zero, $u_{00}(t)$ will be very big; to avoid this, $u_{00}(t)$ is modified as follows:

$$
\begin{equation*}
u_{00}(t)=\frac{-b s}{k_{0} k_{p}\left(s^{2}+\bar{\epsilon}\right)} \sum_{i=1}^{2}\left[\left(\left|\hat{\rho}_{i}\right|+\bar{\rho}_{i}\right)^{2}+\left(\left|\hat{\tau}_{i}\right|+\bar{\tau}_{i}\right)^{2}+\left(\left|\hat{\beta}_{i}\right|+\bar{\beta}_{i}\right)^{2}\right] \tag{46}
\end{equation*}
$$

in which $\bar{\epsilon}$ is a small positive number.
The update laws for delays in Equation (27), which are not available, depend on the estimation error. This problem can be solved by the following:

Given that $0<\underline{\tau_{i}}<\tau_{i}<\tau_{i}$, such that $\tau_{i}$ is the upper limit and $\underline{\tau_{i}}$ is the lower limit of the time delay, as a result of selecting $\hat{\tau}_{i}(0)=\overline{\tau_{i}}$, we have:

$$
\tilde{\tau}_{i}(0)=\tau_{i}-\hat{\tau}_{i}(0)=\tau_{i}-\bar{\tau}<0 \Rightarrow \operatorname{sgn}\left(\tilde{\tau}_{i}\right)=-1
$$

By defining $V_{\tilde{\tau}_{i}}=\frac{1}{2} \tilde{\tau}_{i}^{2}$ and calculating its derivate:

$$
\begin{equation*}
D^{q} V_{\tilde{\tau}_{i}} \leq \tilde{\tau}_{i} D^{q} \tilde{\tau}_{i}=-\tilde{\tau}_{i} l_{i}|s| \operatorname{sgn}\left(\tilde{\tau}_{i}\right)=-l_{i}\left|\tilde{\tau}_{i}\right||s|<0 \tag{47}
\end{equation*}
$$

Therefore, $V_{\widetilde{\tau}_{i}}$ is a decreasing function that tends to zero as a result: $\forall t \geq 0: \tilde{\tau}_{i}<0 \Rightarrow$ $\operatorname{sgn}\left(\tilde{\tau}_{i}\right)=-1$.

In this way, the update laws for time delays are as follows:

$$
\begin{equation*}
D^{q} \hat{\tau}_{i}=l_{i}\left|\operatorname{sisgn}\left(\tilde{\tau}_{i}\right)=-l_{i}\right| s \mid i=1.2 \tag{48}
\end{equation*}
$$

Also, in order to increase the robustness of the adaptive laws against uncertainties and disturbances, the Sigma correction law was used. The behavior of the sigma function is shown in Figure 1.

The sigma function is defined as follows:

$$
\sigma(t)=\left\{\begin{array}{c}
0 \text { if }|\hat{\theta}(t)| \leq M_{0}  \tag{49}\\
\left(|\hat{\theta}(t)| / M_{0}-1\right)^{n} \sigma_{0} \text { if } M_{0}<|\hat{\theta}(t)| \leq 2 M_{0} \\
\sigma_{0} \text { if }|\hat{\theta}(t)| \geq 2 M_{0}
\end{array}\right.
$$



Figure 1. Sigma function behavior diagram.
Therefore, the update laws for estimations of delays, disturbance, and uncertainty bounds are as follows:

$$
\begin{align*}
& D^{q} \hat{\tau}_{i}=-l_{i}|s|-\sigma_{0}\left(\left|\hat{\tau}_{i}\right|\right) \hat{\tau}_{i}, \hat{\tau}_{i}(0)=\overline{\tau_{i}}, i=1.2 \\
& D^{9} \hat{\rho}_{i}=k_{0} k_{p}|s|-\sigma_{0}\left(\left|\hat{\rho}_{i}\right|\right) \hat{\rho}_{i}, i=1.2 \\
& D^{q} \hat{\beta}_{1}=-k_{0} k_{p}|s| \omega_{1}(x)-\sigma_{0}\left(\left|\hat{\beta}_{1}\right|\right) \hat{\beta}_{1},  \tag{50}\\
& D^{q} \hat{\beta}_{2}=-k_{0} k_{p}|s| \omega_{2}(y)-\sigma_{0}\left(\left|\hat{\beta}_{2}\right|\right) \hat{\beta}_{2},
\end{align*}
$$

Its stability is demonstrated for chaotic systems with unknown uncertainty, fractional order unknown time delay, and considering PI sliding surface and nonlinear fractional order derivative.

## 6. Simulation Results

In this section, the process of synchronizing time-varying chaotic systems with unknown uncertainty and time delay of the fractional order using the proposed control mechanism based on the nonlinear fractional order PID and with the advantage of the adaptive controller and update laws that estimate system parameters is verified, and its accuracy is evaluated. Two modified Jerk chaotic systems with the aforementioned characteristics were utilized for this purpose. The canonical form of the master system's governing equations are as follows [15]:

$$
\left\{\begin{array}{c}
D^{q} x_{1}=x_{2}  \tag{51}\\
D^{q} x_{2}=x_{3} \\
D^{q} x_{3}=-\varepsilon_{1} x_{1}(t)-x_{2}(t)-\varepsilon_{2} x_{3}(t)+f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right)
\end{array}\right.
$$

In this system, $f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right)$ is a piecewise linear function, as follows:

$$
\begin{equation*}
f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right)=\frac{1}{2}\left(v_{0}-v_{1}\right)\left[\left|x_{1}\left(t-\tau_{1}\right)+1\right|-\left|x_{1}\left(t-\tau_{1}\right)-1\right|\right]+v_{1} x_{1}\left(t-\tau_{1}\right) \tag{52}
\end{equation*}
$$

Such that $v_{0}<-1<v_{1}<0, v_{0}=-2.5$, and $v_{1}=-0.5$.
Also, $\varepsilon_{i}$ is a time-varying function, defined as follows.

$$
\begin{align*}
& \varepsilon_{1}(t)=0.5+0.3 \sin (t) \cos (5 \pi t)  \tag{53}\\
& \varepsilon_{2}(t)=0.2+0.15 \sin (0.5 t) \cos (3 \pi t)
\end{align*}
$$

If $\varepsilon_{i}(t)=\varepsilon_{0 i}+\Delta \varepsilon_{i}(t), \Delta \varepsilon_{i}(t)$ can be considered a part of the uncertainty and summed with the general uncertainty.

Thus, Equation (49) can be rewritten as follows:

$$
\left\{\begin{array}{c}
D^{q} x_{1}=x_{2}  \tag{54}\\
D^{q} x_{2}=x_{3} \\
D^{q} x_{3}=-\epsilon_{10} x_{1}(t)-x_{2}(t)-\epsilon_{20} x_{3}(t)+f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right) \\
+\Delta f^{n e w}(x(t), t)+d_{1}(t)
\end{array}\right.
$$

in which $\left.\Delta f^{n e w}(x(t), t)\right)=\Delta f(x(t), t)-\Delta \varepsilon_{1}(t) x_{1}(t)-\Delta \varepsilon_{2}(t) x_{3}(t)$ with the previous structure. The same is carried out for the slave system.

When the initial conditions are chosen as $\left(x_{1}(0) ; x_{2}(0) ; x_{3}(0)\right)^{T}=(-0.5032 ; 2.8545$; $-1.37)^{T}$, the chaotic behavior of the system is as shown in Figure 2.


Figure 2. Chaotic behavior of the fractional order Jerk master and slave systems without applying the controller.

If the bounded uncertainty functions of the master and slave systems are as follows:

$$
\begin{align*}
& \Delta f(x(t), t))=0.3 \sin \left(4 x_{1}(t)+x_{2}(t)-x_{3}(t)\right)  \tag{55}\\
& \Delta g(x(t), t))=0.2 \sin \left(y_{1}(t)+2 y_{2}(t)-y_{3}(t)\right)
\end{align*}
$$

The dynamic equations of the master and slave system are as follows:

$$
\left\{\begin{array}{c}
D^{q} x_{1}=x_{2}  \tag{56}\\
D^{q} x_{2}=x_{3} \\
D^{q} x_{3}=-\varepsilon_{1} x_{1}(t)-x_{2}(t)-\varepsilon_{2} x_{3}(t)+f_{3}\left(x_{1}\left(t-\tau_{1}\right), t\right) \\
\left.+\Delta f^{n e w}(x(t), t)\right)+d_{1}(t)
\end{array}\right.
$$

The dynamic of the master system follows the following equations:

$$
\left\{\begin{array}{c}
D^{q} y_{1}=y_{2}  \tag{57}\\
D^{q} y_{2}=y_{3} \\
D^{q} y_{3}=-\varepsilon_{1} y_{1}(t)-y_{2}(t)-\varepsilon_{2} y_{3}(t)+g_{3}\left(y_{1}\left(t-\tau_{2}\right), t\right) \\
\left.+\Delta g^{n e w}(x(t), t)\right)+d_{2}(t)+u(t)
\end{array}\right.
$$

Such that the nonlinear terms of the slave system are as follows:

$$
\begin{equation*}
g_{3}\left(y_{1}\left(t-\tau_{2}\right), t\right)=\frac{1}{2}\left(v_{0}-v_{1}\right)\left[\left|y_{1}\left(t-\tau_{2}\right)+1\right|-\left|y_{1}\left(t-\tau_{2}\right)-1\right|\right]+v_{1} y_{1}\left(t-\tau_{2}\right) \tag{58}
\end{equation*}
$$

According to the dynamic of the master and slave systems described in Equations (51) and (52), the synchronization error is given as follows:

$$
\left\{\begin{array}{c}
D^{q} e_{i}=e_{i+1} 1 \leq i \leq n-1  \tag{59}\\
D^{q} e_{n}=\sigma_{0}^{T} \cdot e(t)+g\left(y\left(t-\tau_{2}\right), t\right)+\Delta g(y(t), t)+d_{2}(t) \\
\left.-f\left(x_{1}\left(t-\tau_{1}\right)\right)-\Delta f^{n e w}(x(t), t)\right)-d_{1}(t)+u(t)
\end{array}\right.
$$

Accordingly, the error dynamics for the chaotic Jerk system are as follows:

$$
\left\{\begin{array}{c}
D^{q} e_{1}=e_{2}  \tag{60}\\
D^{q} e_{2}=e_{3} \\
D^{q} e_{3}=-\varepsilon_{1} e_{1}(t)-e_{2}(t)-\varepsilon_{2} e_{3}(t)-g\left(y_{1}\left(t-\tau_{2}\right)\right)+f\left(x_{1}\left(t-\tau_{1}\right)\right) \\
\left.\left.+\Delta g^{n e w}(x(t), t)\right)-\Delta f^{n e w}(x(t), t)\right) \\
d_{2}(t)-d_{1}(t)+u(t)
\end{array}\right.
$$

At this stage, we applied the robust adaptive control signal, which is devised by combining the sliding surface based on the structure of the fractional order nonlinear PID controllers and described in Equation (26), to the slave system.

In this article, simulations were run for 100 s . Figure 2 depicts the master and slave systems in three-dimensional space. Figure 3 illustrates the behavior of the master and slave system states in the absence of any controller actions. Figure 4 shows the synchronization of the master and slave system. It is clear that after applying the control signal based on the proposed mechanism, the slave system follows the master system well. Figure 5 depicts the synchronization error of the master and slave system utilizing the proposed mechanism. Figure 6 depicts the control signal based on the proposed method. According to the range of the image's control signal (6), it is unquestionable that the proposed controller can be implemented. As this figure demonstrates, the controller signal exhibited no chattering, and a saturation limit of 24 volts was used, which is simple to implement. In this design, the controller coefficients $k_{11}=k_{22}=9$ and $k_{12}=k_{21}=18$ were selected. Also, the gain and time constants of the PID sliding surface are nonlinear fractional orders, as $k_{p}=3$, $T_{i}=0.8$, and $T_{d}=0.65$. The fractional order of the integral part and the derivative of the sliding surface are defined as $\delta=0.15$ and $\lambda=0.75$. The parameters of the proposed robust controller are $\bar{\epsilon}=0.01$ and $b=2$. The unknown time delays of the system are $\tau_{1}=0.3$ and $\tau_{2}=0.5$. The time delay of the master system changes to the value of $\tau_{1}=0.45$ at the moment $t=40 \mathrm{~s}$, and the time delay of the follower system changes to the value of $\tau_{2}=0.58$ at the moment $t=50 \mathrm{~s}$. The error in estimating the uncertainty, disturbance, and delay bounds is shown in Figure 7. Figure 8 shows the uncertainties and disturbances applied to the master and slave systems. The unknown disturbances are applied to both systems as follows:

$$
\begin{align*}
& d_{1}(t)=0.8 \sin ^{2} 2 t+1.2 \cos 3 t+1.6 \sin 1.3 t \\
& d_{2}(t)=\sin 1.4 t+0.3 \sin \pi t+0.3 \cos \frac{\pi t}{2} \tag{61}
\end{align*}
$$



Figure 3. The behavior of the master and slave system states without applying the control signal.


Figure 4. Synchronization of chaotic jerk systems with the help of the proposed control mechanism and application of the control signal at $\mathrm{t}=5 \mathrm{~s}$.


Figure 5. Synchronization error of the master and slave systems using the proposed adaptive sliding mode control mechanism.


Figure 6. Control signal based on the proposed adaptive sliding mode control mechanism.


Figure 7. System parameter estimation error including time delay, disturbance bound, and uncertainty bound.


Figure 8. Uncertainties and disturbances in the master and slave systems.

## 7. Application of Secure Communication in Encryption and Image Retrieval

Despite the uncertainties and time delays in the system, the fractional order chaotic master and slave systems were entirely synchronized according to the proposed mechanism, the details of which were described in the previous section. Images were encrypted using the [27] algorithm in this section. The encrypted image was then transmitted using fractional order chaotic masking and received with high precision before being decoded.

Figure 9 is a block diagram detailing the encryption technique applied to the images. In this block diagram, information is exchanged via a wireless communication channel.

Various statistical parameters, including the histogram difference between the original image and the restored image, correlation, NPCR, PSNR, and information entropy, were calculated for standard color benchmark images and medical color images to demonstrate the efficacy of the proposed method. These parameters are standard criteria that have been used in numerous articles [27].

This section encrypts images for secure communication utilizing the mechanism whose efficacy was evaluated in Section 6. Figure 10 shows the result of image encryption and recovery using secure communication for the original image, and Figure 11 shows their histogram for Aletta (Isekai.Shokudou) color image.


Figure 9. Block diagram of chaotic masking for image encryption.


Figure 10. The original, encrypted, and decrypted color image.


Figure 11. Histogram of the original, encrypted, and decrypted color image.
Figure 12 shows the encryption on the lena color image and Figure 13 shows its histogram.


Figure 12. Original, encrypted, and decrypted color image.


Figure 13. Histogram of the original, encrypted, and decrypted color image.
It can be seen that the decoded images were well restored using the proposed synchronization scheme.

Table 2 shows the results of the statistical criteria of Figures 10 and 12.
Table 2. Results of statistical criteria of color images.

| Images | Histogram |  | Correlation | Differential Attack |  | PSNR | Information Entropy |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Standard | Encrypted |  | NPCR (\%) | UACI (\%) |  |  |
| Images 10 | 21,153.1171 | 21,148.239 | 0.0068 | 99.68 | 33.23 | 8.10 | 7.9690 |
| Images 12 | 18,144.3510 | 18,143.750 | 0.0043 | 99.40 | 33.46 | 8.27 | 7.9700 |

Image encryption using the above mechanism along with histogram for cameraman's black and white image is shown in Figures 14 and 15, respectively, and for panda is presented in Figures 16 and 17.


Figure 14. The original, encrypted, and decrypted black and white image.


Figure 15. Histogram of the original, encrypted, and decrypted black and white image.


Figure 16. The original, encrypted, and decrypted black and white image.


Figure 17. Histogram of the original, encrypted, and decrypted black and white image.
Table 3 shows the results of the statistical criteria of the black and white images in Figures 14 and 16.

Table 3. Results of statistical measures of black and white images.

| Images | Histogram |  | Correlation | Differential Attack |  | PSNR | Information Entropy |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Main | Decoded |  | NPCR (\%) | UACI (\%) |  |  |
| Images 14 | 398,232.09375 | 398,201.1053 | 0.9923 | 99.21 | 33.55 | 8.9671 | 7.9783 |
| Images 16 | 24,466.718750 | 24,421.32934 | 0.9953 | 99.48 | 33.21 | 8.0221 | 7.9458 |

The encryption of the medical color image along with its histogram for the single image mode is shown in Figures 18 and 19 and for the multiple image in Figures 20 and 21, respectively.


Figure 18. Original, encrypted, and decrypted medical image.


Figure 19. Histogram of the original, encrypted, and decrypted medical image.


Figure 20. The original, encrypted, and decrypted medical image.


Figure 21. Histogram of the original, encrypted, and decrypted medical image.
Table 4 shows the results of the statistical criteria of the color medical images in Figures 18 and 20. The results of encryption entropy indicate excellent quality of image retrieval.

Table 4. Results of statistical measures of medical images.

| Images | Histogram |  | Correlation | Differential Attack |  | PSNR | Information Entropy |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Standard | Decrypted |  | NPCR (\%) | UACI (\%) |  |  |
| Images 18 | 65,536 | 65,535 | 0.9986 | 99.96 | 33.46 | 9.23 | 7.9627 |
| Images 20 | 65,536 | 65,534 | 0.9987 | 99.97 | 33.47 | 9.24 | 7.9842 |

## 8. Conclusions

This study examined a novel adaptive sliding mode control approach for robust synchronization of a class of fractional order chaotic systems with uncertainty, external disturbance, and unknown parameters, such as unknown time delay. In the proposed robust control mechanism, a nonlinear fractional order sliding surface was first proposed based on the structure of nonlinear proportional, integral, and fractional derivative controllers. Using the Lyapunov theory and Lipschitz conditions in chaotic systems, matching criteria were established in order to estimate the unknown parameters of the system. In order to facilitate the implementation process, the control signal's saturation limit was defined, and the robust control system's stability was demonstrated. The synchronization of two fractional order Jerk chaotic systems with the stated characteristics, including uncertainties and unknown time delays, based on the proposed control mechanism was simulated using MATLAB, and the results express the capability and optimal performance of the proposed approach in the robust synchronization of the mentioned systems. In closing, the proposed adaptive sliding mode control approach was implemented in the structure of a chaotic secure communication mechanism, and the simulation results indicate a high level of quality in the secure encryption and decryption of digital images despite the presence of uncertain parameters in the master and slave systems of the communication mechanism.


#### Abstract

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## Article

# Analysis and Design of Robust Controller for Polynomial Fractional Differential Systems Using Sum of Squares 

Hassan Yaghoubi ${ }^{1(\mathbb{D}}$, Assef Zare ${ }^{1, *(\mathbb{D})}$ and Roohallah Alizadehsani ${ }^{2(D)}$<br>1 Department of Electrical and Electronic Engineering, Gonabad Branch, Islamic Azad University, Gonabad 6518115743, Iran<br>2 Institute for Intelligent Systems Research and Innovation (IISRI), Deakin University, Victoria 3216, Australia * Correspondence: assefzare@gmail.com

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#### Abstract

This paper discusses the robust stability and stabilization of polynomial fractional differential (PFD) systems with a Caputo derivative using the sum of squares. In addition, it presents a novel method of stability and stabilization for PFD systems. It demonstrates the feasibility of designing problems that cannot be represented in LMIs (linear matrix inequalities). First, sufficient conditions of stability are expressed for the PFD equation system. Based on the results, the fractional differential system is Mittag-Leffler stable when there is a polynomial function to satisfy the inequality conditions. These functions are obtained from the sum of the square (SOS) approach. The result presents a valuable method to select the Lyapunov function for the stability of PFD systems. Then, robust Mittag-Leffler stability conditions were able to demonstrate better convergence performance compared to asymptotic stabilization and a robust controller design for a PFD equation system with unknown system parameters, and design performance based on a polynomial state feedback controller for PFD-controlled systems. Finally, simulation results indicate the effectiveness of the proposed theorems.


Keywords: polynomial fractional-order system; robust controller; stability; stabilization; Mittag-Leffler stable

MSC: 93D09; 93B51

## 1. Introduction

Fractional calculus concerns mathematical relations about the generalizations of differentiation and integration to a noninteger order with a history of more than 300 years. Integer-order derivatives and integrals as specific cases paved the way for this mathematical branch to become very popular in fractional calculus, which resulted in many applications in engineering, physics, economics, etc. [1]. In addition, new possibilities have caused fractional calculus to model various physical systems in engineering, which have more accuracy than the classical integer system such as a robot, chaos, information science, and so on $[2,3]$. Moreover, new methods were proposed to solve the complexity of modeling by the fractional-order method [4]. Recently, the study of the stability and stabilization of fractional differential equations (FDEs) has attracted a lot of attention in control theory [5,6]. To this aim, many studies have focused on linear fractional differential and nonlinear fractional differential equations [7], which could conform to linear FDE systems and analyze stability based on LMI conditions.

Zhang, Tian, et al. [8] considered the stability of nonlinear FDEs, and similar stability conditions of Caputo FDEs were obtained for Riemann-Liouville FDEs. Based on the result, the stability condition of nonlinear FDEs is the same as linear FDEs if the nonlinear section follows the same conditions. Wang et al. [9] investigated the asymptotical stability of nonlinear FDEs and sufficient conditions obtained by using the state feedback stabilization
controller. Furthermore, the pole replacement method was used in linear systems to design the controller gains for nonlinear FDEs [10].

Nowadays, stability analysis of nonlinear FDE systems has been considered by researchers. Thus, most of the studies are related to stability and stabilization for the fractional differential equation [11]. FDE is analyzed by Lyapunov's first and second methods. In the first method, nonlinear FDEs are converted to linear FDEs at the equilibrium point. Therefore, the nonlinear FDE is asymptotically stable if the linearization system is asymptotically stable $[12,13]$. In the second method, energy is decreased and allows us to evaluate the stability of the system without integrating the differential equation explicitly. The Lyapunov technique provides a sufficient condition for the asymptotic stabilization of systems. In this regard, the LMI approach can be used as a method for selecting a Lyapunov candidate. The LMI method is based on numerical solutions and optimization due to its popularity. Various studies have been conducted in the field of stability evaluation by the LMI method. In the study of Lu and Chen [14], less conservative conditions have been evaluated in terms of LMIs for robust stability and stabilization of FO dynamic interval systems. Furthermore, Li and Zhang [15] presented robust stability of the FO linear uncertain system by focusing on the observer and obtaining the necessary conditions. All of the results in some studies [14-16] were obtained based on LMI, although many design problems cannot be represented by the LMI approach. The analysis stability by using Lyapunov's method is considered an explicit way of solving the FDE in nonlinear FDE systems. Thus, Mittag-Leffler introduced stability for nonlinear fractional differential systems by the fractional Lyapunov's method [17]. Based on this method, systems are stable but have no candidate for the Lyapunov function [18]. Two theorems were proved for fractional nonlinear time-delay systems that were related to stability in the study of Badri and Tavazoei [19]. M-L stability for nonlinear FDE systems can generalize better convergence performance against asymptotic stabilization. Chen et al. [20] studied the stabilization of fractional nonautonomous systems by using the M-L function and the Lyapunov direct method. Some studies reported that it is usually difficult to find a Lyapunov candidate and calculate a fractional derivative for the FDE system (e.g., [10,18-20]). By considering all of the above-mentioned studies, a new method was presented for finding a Lyapunov candidate function. This method can help find the Lyapunov function more easily than the previous methods. The result is based on the new property of the Caputo fractional derivative, which allows the stability analysis of many FO systems to be studied [21,22].

Some studies reported that stability analysis based on M-L stability is more efficient than asymptotic stabilization [17,19-23]. In this paper, the application of the Lyapunov function method was expanded in PFD systems. To this aim, the stability of the PFD system was analyzed by using three polynomial PFD inequalities, which can be solved via the SOS toolbox in Matlab. Then, robust Mittag-Leffler stability conditions were obtained based on the SOS approach, which can exhibit better convergence. The desired robust M-L stabilization was obtained by selecting the polynomials state feedback control, which resulted in designing flexible controllers.

This paper is organized as follows. In Section 2, some definitions and lemmas are given. A sufficient condition of stability for PFD is given in Section 3. Section 4 provides sufficient conditions for robust stability and the stabilization of the PFD system. Simulation results are given in Section 5. Finally, some conclusions are made in Section 6.

## 2. Preliminaries

Notations and Definitions
There are several definitions of FO derivatives, among which Riemann-Liouville and Caputo's definition is considered the most common and practical definition in the literature. Thus, the Caputo definition is selected in this study.

Definition 1 ([8]). The Caputo fractional derivative is defined as follows

$$
\begin{equation*}
D_{t}^{q} f(t)=\frac{1}{\Gamma(n-q)} \int_{t_{0}}^{t}(t-\tau)^{n-q-1} f^{(n)}(\tau) d \tau \tag{1}
\end{equation*}
$$

where $q$ is the fractional order and $n$ shows integer.

$$
n-1 \leq q<n
$$

$\Gamma(\cdot)$ represents the Euler's function

$$
\begin{equation*}
\Gamma(t)=\int_{0}^{\infty} x^{t-1} e^{-x} d x \tag{2}
\end{equation*}
$$

Mittag-Leffler is a function that is mostly used in solving in fractional-order systems as follows

$$
\begin{equation*}
E_{q}(z)=\sum_{k=0}^{\infty} \frac{z^{k}}{\Gamma(k q+1)} . \tag{3}
\end{equation*}
$$

where $q>0$. The M-L function with two parameters appears most frequently and has the following form

$$
\begin{equation*}
E_{q, \beta}(z)=\sum_{k=0}^{\infty} \frac{z^{k}}{\Gamma(k q+\beta)} . \tag{4}
\end{equation*}
$$

where $q>0, \beta>0$.
By using the Caputo derivative, an FO system is defined by

$$
\begin{equation*}
D_{t}^{q} x(t)=(A(x)+\Delta A(x)) x(t)+u(x) \tag{5}
\end{equation*}
$$

where $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{T} \in R^{n}$ is the state vector of the state system, $A(x) \in R^{n \times n}$ defines a nonlinear matrix function field in the $n \times n$-dimensional space, and $\Delta A(x)$ represents the admissible uncertainty function. $q$ is the order of the fractional derivative, $(0<q \leq 1) u(x)$ as the control input. Let the equilibrium point be $x=0$ when $u(x)=0$.

Definition 2 ([24]). The sum of squares (SOS) approach is an important subset of the polynomials used for modeling and controlling nonlinear systems. Assume that $\sum_{n}$ is the set of all SOS polynomials with degree $n$ defined as follows

$$
\sum_{n}=\left\{s \in R_{n} \mid \exists M \prec \infty, \exists\left\{p_{i}\right\}_{i=1}^{M} \subset R_{n} \text { such that } s=\sum_{i=1}^{M} p_{i}^{2}\right\}
$$

where $R$ indicates the real number. Assume that monomial $\left\{m_{\alpha_{j}}\right\}_{j=1}^{k}$ is defined as $m_{\alpha}(x)=x^{\alpha}=$ $x_{1}^{\alpha_{1}} . x_{2}^{\alpha_{2}} \ldots . x_{n}^{\alpha_{n}}, \alpha \in Z_{+}$. Then, polynomial $p$ is the issue of squares if $p$ is a monomial $\left\{m_{\alpha_{j}}\right\}_{j=1}^{k}$ linear combination.

$$
\begin{aligned}
& \left\{m_{\alpha_{j}}\right\}_{j=1}^{k} \text { so that } m_{\alpha}: R^{n} \rightarrow R \text { subject to } \\
& \qquad \begin{array}{l}
m_{\alpha}(x)=x^{\alpha}=x_{1}^{\alpha_{1}} \cdot x_{2}^{\alpha_{2}} \ldots x_{n}^{\alpha_{n}}, \alpha \in Z_{+},\left\{c_{j}\right\}_{j=1}^{k} \in R, \\
p=\sum_{j=1}^{k} c_{j} m_{\alpha_{j}} \\
p=\sum_{j=1}^{k} c_{j} m_{\alpha_{j}},\left\{c_{j}\right\}_{j=1}^{k} \in R
\end{array}
\end{aligned}
$$

A subset of $R_{n}, R_{n, d}=\left\{p \in R_{n} \mid \operatorname{deg} p \leq d\right\}$ is defined in such a way that $n$ is the number of variables and d is the degree of polynomials. Based on this definition, we can
define $\sum_{n, d}$, which can directly lead to sufficient conditions for polynomial programming. Therefore, we have $\sum_{n, d}=\sum_{n} \cap R_{n, d}$.

Lemma 1 ([11]). Fix $p \in R_{n, 2 d} . p \in \sum_{n, 2 d}$ if and only if there exists a $Q \geq 0$ such that

$$
\begin{equation*}
p(x)=z^{*}{ }_{n, d}(x) Q z_{n, d}(x) \tag{6}
\end{equation*}
$$

According to Lemma 1, polynomial $p$ is SOS if necessary and sufficient conditions are satisfied for Lemma 1.
$z_{n, d}(x)$ is monomials with $n$ variables of degree less than or equal to $d$.
Theorem 1 ([11]). The polynomial system $\dot{x}(t)=f(x(t))$ is globally asymptotically stable about equilibrium point if there exists a positive-definite function $V: R^{n} \rightarrow R_{+}$such that $-\dot{v}(x)$ is positive-definite.

This theorem is important in stability, which is defined based on the following important theorem in the field of stability of polynomial systems.

Theorem 2. Given the system $\dot{x}(t)=f(x(t))$ and fixed positive-definite functions $l_{1}, l_{2} \in R^{n}$, the system is globally asymptotically stable if there exists $v(x) \in R^{n}$ with $v(0)=0$ such that

$$
\begin{gather*}
v(x)-l_{1} \in \sum_{n} . \\
-\left(\dot{v}(x)+l_{2}\right) \in \sum_{n} .  \tag{7}\\
\dot{v}(x)=\Delta v(x) A(x) x
\end{gather*}
$$

Proof. Given a finite set $\left\{p_{i}\right\}_{i=0}^{m} \in R_{n}$, the existence of $\left(p_{0}+\sum_{i=1}^{m} \alpha_{i} p_{i}\right) \in \sum_{n}$ is such that $\left\{\alpha_{i}\right\}_{i=1}^{m} \in R$. It is evident that the conditions $\left(v(x)-l_{1}\right)$ and $-\left(\nabla v(x) A(x)+l_{2}\right)$ are SOS polynomials if a polynomial function of $v(x)$ is found for satisfying these conditions. The positive definiteness of $l_{1}$ and $l_{2}$ is selected to satisfy the assumptions of theorem, in which both $v(x)$ and $-\dot{v}(x)$ are positive-definite.

Definition 3 ([20]). Assume that $v(x): \Omega \rightarrow R, \Omega \in R^{n}$ is a convex function in $\Omega$ and $x:\left[t_{0}, \infty\right) \rightarrow \Omega$ is a continuous differential function. For $t \geq t_{0}$

$$
\begin{equation*}
{ }_{t_{0}}^{c} D_{t}^{q} v(x) \leq\left(\frac{\partial v}{\partial x}\right)^{T}{ }_{t_{0}}^{c} D_{t}^{q} x(t) \tag{8}
\end{equation*}
$$

Definition 4 ([20]). Consider the FO nonlinear system

$$
\begin{equation*}
{ }_{t_{0}}^{c} D_{t}^{q} x(t)=A(x) x(t) . \tag{9}
\end{equation*}
$$

If the convex Lyapunov function $v(x)$ can satisfy the following conditions, $x=0$ is an equilibrium point for the FO system and the system is globally Mittag-Leffler stable in equilibrium point.

$$
\begin{gather*}
\gamma_{1}\|x\|^{\alpha} \leq v(x) \leq \gamma_{2}\|x\|^{\alpha b} .  \tag{10}\\
c_{t_{0}} D_{t}^{q} v(x) \leq-\gamma_{3}\|x\|^{\alpha b} .
\end{gather*}
$$

## 3. PFDE Stability

In this section, a sufficient condition of stability is presented for the PFD equations system.

Theorem 3. Consider the system $c_{t_{0}} D_{t}^{q} x(t)=A(x) x(t)$. Suppose that Let $x=0$ is an equilibrium point in the domain $D \subset R^{n}$ and $v(x):[0, \infty) \times D \rightarrow R$ is a continuously differentiable
function and locally Lipschitz in $x$. System is $M-L$ stable in $x=0$ if and only if there exists $w_{1}(x), w_{2}(x), w_{3}(x) \in \sum_{n}$ by satisfying the following condition

$$
\begin{align*}
& w_{1}(x) \leq v(x) \leq w_{2}(x) \\
& c_{t_{0}} D_{t}^{q} v(x) \leq-w_{3}(x) \tag{11}
\end{align*}
$$

Proof. It follows from inequalities (11) that

$$
\begin{equation*}
{ }_{t_{0}}^{c} D_{t}^{q} v(x) \leq-w_{3}(x) \leq \frac{-w_{3}(x)}{w_{2}(x)} v(x) \leq-i n f_{x}\left(\frac{w_{3}(x)}{w_{2}(x)}\right) v(x)=-\beta v(x) . \tag{12}
\end{equation*}
$$

where $\beta=\inf _{x}\left(\frac{W_{3}(x)}{W_{2}(x)}\right)$.
$w_{2}(x), w_{3}(x)$ are positive and have the same degree, therefore: $\beta>0$.
There exists a nonnegative function $m_{0}(t)$ satisfying

$$
\begin{equation*}
{ }_{t_{0}}^{c} D_{t}^{q} v(x)+m_{0}(t)=-\beta v(x) . \tag{13}
\end{equation*}
$$

By taking the Laplace transform from both sides of Equation (13)

$$
\begin{gather*}
v(s)=\frac{v(0) s^{q-1}-M_{0}(s)}{s^{q}+\beta}  \tag{14}\\
s^{q} v(s)-v(0) s^{q-1}+M(s)=-\beta v(s)
\end{gather*}
$$

Then, if $x(0)=0$, then $v(0)=0$, and if $x \neq 0$, then $v(0)>0$. Applying the inverse Laplace transform to (14) and according to Definition 1 gives

$$
\begin{equation*}
v(t)=v(0) E_{q}\left(-\beta t^{q}\right)-M(t)\left[t^{q-1} E_{q}\left(-\beta t^{q}\right)\right] . \tag{15}
\end{equation*}
$$

$t^{q-1} \geq 0, E_{q} \geq 0$ are nonnegative functions. It follows that

$$
\begin{equation*}
v(t) \leq v(0) E_{q}\left(-\beta t^{q}\right) . \tag{16}
\end{equation*}
$$

Substituting (16) into (11) yields

$$
\begin{equation*}
w_{1}(x) \leq v(0) E_{q}\left(-\beta t^{q}\right)=v(0) \sum_{k=0}^{\infty} \frac{\left(-\beta t^{q}\right)^{k}}{\Gamma(k q+1)} \tag{17}
\end{equation*}
$$

where $v(0)>0$ for $x(0) \neq 0$.
Therefore $w_{1}(x)$ is bounded, also $v(x)$ is locally Lipschitz in $x . v(0, x(0))=0$ if and only if $x(0)=0$, which guarantees the Mittag-Leffler stability of system (9).

## 4. Robust Stability

In this section, a sufficient condition of robust stability is presented for the PFD equations system. The uncertain FDE system (18) is robust M-L stable if there is SOS polynomial $w_{1}(x), w_{2}(x), w_{3}(x)$ so that conditions (20) are satisfied.

Theorem 4. Consider the FO nonlinear system

$$
\begin{equation*}
{ }_{t_{0}}^{c} D_{t}^{q} x(t)=(A(x)+\Delta A(x, t)) x(t) . \tag{18}
\end{equation*}
$$

where $\Delta A(x, t)$ represents the admissible uncertainty function and where the following condition is satisfied

$$
\begin{equation*}
\|\Delta A(x, t)\| \leq \gamma . \tag{19}
\end{equation*}
$$

This polynomial FO system is robust Mittag-Leffler stability if and only if this condition is satisfied

$$
\begin{gather*}
w_{1}(x) \leq v(x) \leq w_{2}(x)  \tag{20}\\
\left(\frac{\partial v(x)}{\partial x}\right)^{T} A(x) x(t) \leq-w_{3}(x)-\sum_{k=1}^{M} \alpha_{k}\|x\|^{2 k}  \tag{21}\\
\left\|\frac{\partial v(x)}{\partial x}\right\| \gamma\|x\| \leq w(x)  \tag{22}\\
w(x)-\sum_{k=1}^{M} \alpha_{k}\|x\|^{2 k} \leq(1-\eta) w_{3}(x) \tag{23}
\end{gather*}
$$

where: $w_{1}(x), w_{2}(x), w_{3}(x)$ are SOS, $\alpha_{k}>0, M \in \mathbb{N}$ and $\eta \in(0,1)$.
Proof. The convex function $v(x)=x(t)^{T} p(x) x(t)$ is selected as the Lyapunov function for the PFD system $t_{0}^{c} D_{t}^{q} x(t)=(A(x)+\Delta A(x, t)) x(t)$. Based on Definition 3, we have

$$
\begin{gather*}
{\stackrel{c}{t_{0}}}^{D_{t}^{q} v(x) \leq}\left(\frac{\partial v(x)}{\partial x}\right)^{T}(A(x)+\Delta A(x)) x(t)= \\
\left(\frac{\partial v(x)}{\partial x}\right)^{T} A(x) x(t)+\left(\frac{\partial v(x)}{\partial x}\right)^{T} \Delta A(x) x(t) \leq\left(\frac{\partial v(x)}{\partial x}\right)^{T} A(x) x(t)+\left\|\frac{\partial v(x)}{\partial x}\right\| \gamma\|x\|  \tag{24}\\
\leq-w_{3}(x)-\sum_{k=1}^{M} \alpha_{k} x^{2 k}+w(x) \leq-\eta w_{3}(x) .
\end{gather*}
$$

Therefore, according to Theorem 3, the system is Mittag-Leffler stable.
For simplicity, we assume that

$$
\begin{equation*}
h(x)=(1-\eta) w_{3}(x)+\sum_{k=1}^{M} \alpha_{k}\|x\|^{2 k}-w(x) . \tag{25}
\end{equation*}
$$

Therefore, if conditions (20)-(23) are satisfied and $h(x) \geq 0$, the system is MittagLeffler stable. In this theorem, if $\Delta \mathrm{A}(x, t)=0$ then $\gamma=0, \alpha=0, \eta=0$ and conditions of Theorem 4 are converted to Theorem 3 .

## 5. Robust Stabilization

In this section, robust stability PFD equations are studied by designing a polynomial feedback controller.

Theorem 5. Consider the PFDE systems

$$
\begin{equation*}
{ }_{t_{0}}^{c} D_{t}^{q} x(t)=(A(x)+\Delta A(x)) x(t)+u(x) . \tag{26}
\end{equation*}
$$

System (26) is robust Mittag-Leffler stability if and only if the condition of Theorem 4 is satisfied for $\bar{A}(x)=A(x)-k(x)$. By using state feedback controller, $u(x)=-k(x) x, k(x) \in$ $R^{n \times n}$, the closed-loop system including (23) becomes as follows

$$
\begin{equation*}
{ }_{t_{0}}^{c} D_{t}^{q} x(t)=(A(x)-k(x)) x+\Delta A(x) x=\bar{A}(x) x+\Delta A(x) x . \tag{27}
\end{equation*}
$$

The main purpose is designing the controller, which ensures asymptotic stability.
The flowchart of the proposed method is drawn in Figure 1.


Figure 1. The proposed method.

## 6. Examples and Simulations

In this section, by giving examples, we show the efficiency of the methods expressed in Theorems 3-5.

Example 1. Consider the following PFD equations system

$$
\begin{gather*}
c \\
t_{0}^{c} D_{t}^{0.8} x_{1}(t)=-x_{1}-2 x_{1} x_{2}{ }^{2}-2 x_{2}{ }^{3} .  \tag{28}\\
{ }_{t_{0}}^{c} D_{t}^{0.8} x_{2}(t)=-2 x_{1}{ }^{2} x_{2}-2 x_{2}+0.33 x_{2}{ }^{3} . \\
A(x)=\left[\begin{array}{cc}
-1 & -2 x_{1} x_{2}-2 x_{2}{ }^{2} \\
-2 x_{1} x_{2} & 0.33 x_{2}^{2}-2
\end{array}\right], \Delta A(x)=0
\end{gather*}
$$

By using SOSTOOLS we have

$$
\begin{gather*}
v(x)=0.40401 x_{1}^{4}+0.375 x_{1}^{2} x_{2}^{2}+0.866 x_{1}^{2}+0.396 x_{2}^{3} x_{1}-0.026 x_{1} x_{2} \\
+  \tag{29}\\
+0.473 x_{2}^{4}+0.918 x_{2}^{2} .
\end{gather*}
$$

The Lyapunov function obtained is degrees four in this example, whereas no quadratic function can be found for the system.

$$
v(x)=z^{T} Q z \text {, where } z^{T}=\left[\begin{array}{lllll}
x_{1} & x_{2} & x_{1}^{2} & x_{1} x_{2} & x_{2}^{2}
\end{array}\right] \text { and }
$$

$$
Q=\left[\begin{array}{ccccc}
0.8662 & -0.0134 & 0.0000 & 0.0000 & 0.0000  \tag{30}\\
-0.0134 & 0.9180 & 0.0000 & 0.0000 & 0.0000 \\
0.0000 & 0.0000 & 0.4016 & 0.0000 & -0.1046 \\
0.0000 & 0.0000 & 0.000 & 0.5843 & 0.1981 \\
0.0000 & 0.0000 & -0.1046 & 0.1981 & 0.4737
\end{array}\right]
$$

According to Theorem 3:

$$
\begin{align*}
& w_{1}(x)=0.246 x_{1}^{4}+0.22875 x_{1}^{2} x_{2}^{2}+0.571 x_{1}^{2}+0.1702 x_{2}^{3} x_{1}-0.026 x_{1} x_{2}  \tag{31}\\
& +0.010 x_{2}{ }^{4}+0.605 x_{2}{ }^{2} \text {. } \\
& w_{2}(x)=0.741 x_{1}{ }^{8}+0.695 x_{1}{ }^{6} x_{2}{ }^{2}+1.866 x_{1}{ }^{6}+0.3315 x_{2}{ }^{3} x_{1}{ }^{5}-0.0228 x_{1}{ }^{5} x_{2} \\
& +1.575 x_{2}{ }^{4} x_{1}^{4}+1.683 x_{2}^{2} x_{1}^{4}+0.695 x_{1}{ }^{4} x_{2}{ }^{6}+1.866 x_{1}{ }^{2} x_{2}{ }^{4}  \tag{32}\\
& +0.3315 x_{2}{ }^{7} x_{1}-0.0228 x_{2}{ }^{5} x_{1}+0.875 x_{2}{ }^{8}+1.683 x_{2}{ }^{6} \text {. } \\
& w_{3}(x)=-4.77 x_{1}^{4} x_{2}^{2}-1.606 x_{1}^{4}-5.590 x_{1}^{3} x_{2}^{3}+0.053 x_{1}^{3} x_{2}-0.0075 x_{1}^{2} x_{2}^{6} \\
& -5.064 x_{1}^{2} x_{2}^{4}-9.38 x_{1}^{2} x_{2}^{2}-1.73 x_{1}^{2}-0.011 x_{1} x_{2}^{7}  \tag{33}\\
& -1.935 x_{1} x_{2}^{5}-6.192 x_{1} x_{1} x_{2}^{3}+0.0803 x_{1} x_{1} x_{2}^{2}-0.018 x_{2}^{8} \text {. }
\end{align*}
$$

Then $x=0$ is Mittag-Leffler stable.
Figure 2 is drawn for different initial conditions that show the stability of the system (25)


Figure 2. Phase portrait of $x_{1}(t), x_{2}(t)$ for PFDE nonlinear system (28).
Example 2. Consider the following PFD system

$$
\begin{gather*}
{ }_{t_{0}}^{c} D_{t}^{0.8} x_{1}(t)=-4 x_{1}+2 x_{1} x_{2}^{2}+0.2 x_{1} \sin \left(x_{1}+2 x_{2}\right) . \\
t_{0} D_{t}^{0.8} x_{2}(t)=-x_{2}-2 x_{1}^{2} x_{2}+0.3 x_{2} \cos x_{1} . \tag{34}
\end{gather*}
$$

Thus $A(x)=\left[\begin{array}{cc}-4 & 2 x_{1} x_{2} \\ -2 x_{1}^{2} x_{2} & -1\end{array}\right], \Delta A(x(t))=\left[\begin{array}{cc}0.2 \sin \left(x_{1}+2 x_{2}\right) & 0 \\ 0 & 0.3 \cos x_{1}\end{array}\right]$ is the uncertainty function and $\|\Delta A(x)\| \leq 0.36$. The state response of the system (34) with $q$ $=0.8$ demonstrate the instability of the system. According to Equation (26), we obtain polynomial controller $u_{1}(x)=2.98 x_{1}+4.0761 x_{2}^{2} x_{1}$ and obtain the condition of Theorem 4 by using SOSTOOLS.

The degree of Lyapunov function is 6 by using $\eta=0.01, \alpha_{1}=3.6$ in Equations (21) and (23).

$$
\begin{gather*}
v(x)=0.0328 x_{1}^{6}+0.3831 x_{1}^{4} x_{2}^{2}+0.3287 x_{1}^{4}+1.211 x_{1}^{2} x_{2}^{4} 1.512 x_{1}^{2} x_{2}^{2} \\
+0.9734 x_{1}^{2}+1.285 x_{2}^{6}+1.313 x_{2}^{4}+1.276 x_{2}^{2} .  \tag{35}\\
w_{1}(x)=0.001 x_{1}^{6}+0.1731 x_{1}^{4} x_{2}^{2}+0.0187 x_{1}^{4}+0.134 x_{1}^{2} x_{2}^{4}+0.414 x_{1}^{2} x_{2}^{2}  \tag{36}\\
+0.253 x_{1}^{2}+0.175 x_{2}^{6}+0.115 x_{2}^{4}+0.575 x_{2}^{2} .
\end{gather*}
$$

$$
\begin{gather*}
w_{2}(x)=1.234 x_{1}^{6}+2.485 x_{1}^{4} x_{2}^{2}+1.52 x_{1}^{4}+4.314 x_{1}^{2} x_{2}^{4}+3.812 x_{1}^{2} x_{2}^{2} \\
+0.973 x_{1}^{2}+1.285 x_{2}^{6}+1.313 x_{2}^{4}+3.54 x_{2}^{2} .  \tag{37}\\
w_{3}(x)=0.24 x_{1}^{4} x_{2}^{2}+0.206 x_{1}^{4}+0.744 x_{1}^{2} x_{2}^{4}+1.01 x_{1}^{2} x_{2}^{2}+2.62 x_{1}^{2}+0.39 x_{2}^{4}  \tag{38}\\
+
\end{gather*}
$$

$w_{1}(x), w_{2}(x), w_{3}(x)$ are SOS and $h(x) \geq 0$. Then the closed-loop system (34) is robust M-L stable.

Figure 3 is drawn for different initial conditions that show the stability of the closedloop system (34).


Figure 3. Phase portrait of $x_{1}(t), x_{2}(t)$ for PFDE nonlinear system (34).
Figures 4 and 5 demonstrate that Equation (25) is valid and its value is also valid for different values of $x_{1}(t), x_{2}(t)$.


Figure 4. Surface $h(x)$ for large range of $x_{1}(t), x_{2}(t)$ and check the condition $h(x)>0$ for stability of nonlinear system (34).


Figure 5. Surface $h(x)$ for small range of $x_{1}(t), x_{2}(t)$ and check the condition $h(x)>0$ for stability of nonlinear system (34).

Example 3. Consider the following PFD equations system

$$
\begin{gather*}
{ }_{t_{0}} D_{t}^{0.9} x_{1}(t)=-2 x_{1}-4 x_{1} x_{2}^{2}+0.3 x_{1} \sin (t) \sin \left(x_{1}-x_{2}\right) .  \tag{39}\\
t_{0} D_{t}^{0.9} x_{2}(t)=-x_{2}-2 x_{1}^{2} x_{2}+0.4 x_{2} \cos x_{1} \cos (t)
\end{gather*}
$$

where $\Delta A(x(t), t)=\left[\begin{array}{cc}0.3 \sin (t) \sin \left(x_{1}-x_{2}\right) & 0 \\ 0 & 0.4 \cos (t) \cos x_{1}\end{array}\right]$ is the uncertainty function and $\|\Delta A(x, t)\| \leq 0.5$.

The state response of the system (39) with $q=0.9$ demonstrates the instability of the system.

According to Equation (26), we obtain a polynomial controller

$$
\begin{equation*}
u_{1}(x)=-1.85 x_{1}+7.24 x_{1} x_{2}^{2} . \tag{40}
\end{equation*}
$$

By obtaining $\eta=0.98, \alpha_{1}=0.2$, then

$$
\begin{gather*}
w_{1}(x)=0.002 x_{1}^{6}+0.218 x_{1}^{4} x_{2}^{2}+0.224 x_{1}^{4}+0.325 x_{1}^{2} x_{2}^{4}+0.458 x_{1}^{2} x_{2}^{2} \\
+  \tag{41}\\
+0.281 x_{1}^{2}+0.021 x_{1}^{6}+0.321 x_{1}^{4} x_{2}^{2}+0.236 x_{1}^{4}+0.352 x_{1}^{2} x_{2}^{4} \\
+ \\
w_{2}(x)=3.432 x_{2}^{2}+0.621 x_{1}^{2}+0.321 x_{2}^{6}+0.123 x_{2}^{4}+0.218 x_{2}^{2} \tag{42}
\end{gather*}
$$

The degree of Lyapunov function is 6 by using the conditions of Theorem 4.

$$
\begin{align*}
v(x)=0.007798 x_{1}^{6} & +0.479 x_{1}^{4} x_{2}^{2}+0.4365 x_{1}^{4}+0.8522 x_{1}^{2} x_{2}^{4}+0.6879 x_{1}^{2} x_{2}^{2} \\
& +0.741 x_{1}^{2}+0.07798 x_{1}^{6}+0.479 x_{1}^{4} x_{2}^{2}+0.4365 x_{1}^{4} \\
& +0.8522 x_{1}^{2} x_{2}^{4}+0.67879 x_{1}^{2} x_{2}^{2}+0.7 x_{1}^{2}+0.5874 x_{2}^{6}  \tag{44}\\
& +0.2236 x_{2}^{4}+1.118 x_{2}^{2} .
\end{align*}
$$

Figure 6 is drawn for different initial conditions that show the stability of the closedloop system (39).


Figure 6. Phase portrait of $x_{1}(t), x_{2}(t)$ for nonlinear time-variant system (39).
Figures 7 and 8 demonstrate that Equation (25) is valid and its value is also valid for different values of $x_{1}(t), x_{2}(t)$.


Figure 7. Surface $h(x)$ for small range of $x_{1}(t), x_{2}(t)$ and check the condition $h(x)>0$ for stability of nonlinear time-variant system (31).


Figure 8. Surface $h(x)$ for small range of $x_{1}(t), x_{2}(t)$ and check the condition $h(x)>0$ for stability nonlinear time-variant system (39).

## 7. Advantages of the Proposed Approach and Suggestions for Future Research

In this paper, the method of PFD systems was highlighted. The main advantages of the proposed method are as follows.

First, the convex Lyapunov methods in the literature such as [19,21] are complicated and are not able to handle the stability and control performances. Second, although the results of some studies (e.g., $[9,14,15]$ ) are obtained based on the LMI approach, many design problems cannot be represented by the LMI inequality. This paper focused on the sum of squares approach for finding a Lyapunov candidate function for PFD systems with uncertainty.

Finally, compared with the existing studies on the stability analysis of nonlinear fractional differential systems (e.g., [14]), the results use the Lyapunov quadratic function in the stability analysis. In this paper, higher-order Lyapunov functions were used in the stability analysis of FDE systems based on the sum of squares approach. Following this, by working on PFD fuzzy systems, the stability analysis and stabilization of these systems can be obtained.

## 8. Conclusions

Finding a Lyapunov candidate function is difficult considering the previous methods. In this paper, the stability and stabilization of PFD systems are presently based on the Lyapunov candidate function. This method can provide the Lyapunov candidate function with a degree higher than two. Therefore, they cannot be solved with LMI techniques. Accordingly, the SOS method is used for solving. For this reason, the system (7) is MittagLeffler stable if the Lyapunov functional can satisfy the conditions of Theorem 3. Following this, the sufficient condition of a robust stability system (17) is expressed in Theorem (4). Finally, the polynomial Lyapunov approach is proposed using SOSTOOLS for fractional systems. This method can provide an efficient way for analyzing the stability of an uncertain polynomial fractional system.

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## Article

# On Mond-Weir-Type Robust Duality for a Class of Uncertain Fractional Optimization Problems 

Xiaole Guo ${ }^{1,2}$

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1 School of Mathematics and Statistics, Chongqing Technology and Business University, Chongqing 400067, China; guoxiaole@ctbu.edu.cn
2 Chongqing Key Laboratory of Social Economy and Applied Statistics, Chongqing Technology and Business University, Chongqing 400067, China


#### Abstract

This article is focused on the investigation of Mond-Weir-type robust duality for a class of semi-infinite multi-objective fractional optimization with uncertainty in the constraint functions. We first establish a Mond-Weir-type robust dual problem for this fractional optimization problem. Then, by combining a new robust-type subdifferential constraint qualification condition and a generalized convex-inclusion assumption, we present robust $\varepsilon$-quasi-weak and strong duality properties between this uncertain fractional optimization and its uncertain Mond-Weir-type robust dual problem. Moreover, we also investigate robust $\varepsilon$-quasi converse-like duality properties between them.


Keywords: fractional optimization; robust duality; constraint qualification condition

MSC: 90C29; 90C46

## 1. Introduction

Let $T$ be a nonempty infinite index set. Suppose that $f_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}, i=1, \ldots, p$, and $h_{t}: \mathbb{R}^{n} \rightarrow \mathbb{R}, t \in T$. Let us consider the semi-infinite optimization problem:
(MP)

$$
\left\{\begin{array}{c}
\operatorname{Min}_{\mathbb{R}_{+}^{p}}\left(f_{1}(x), \ldots, f_{p}(x)\right) \\
\text { s.t. } h_{t}(x) \leq 0, \forall t \in T \\
x \in \mathbb{R}^{n}
\end{array}\right.
$$

The study of optimization problem (MP) is a very interesting topic and has been considered extensively by many scholars from different points of view, see [1-13]. However, most semi-infinite optimization models of real-world problems are contaminated by prediction errors or asymmetry knowledge. Thus, it is necessary to consider semiinfinite optimization problems under uncertain data. This optimization problem (MP) with uncertainty can be captured by
(UMP)

$$
\begin{cases}\operatorname{Min}_{\mathbb{R}_{+}^{p}}\left(f_{1}(x), \ldots, f_{p}(x)\right) \\ \text { s.t. } & h_{t}\left(x, v_{t}\right) \leq 0, \forall t \in T \\ & x \in \mathbb{R}^{n}\end{cases}
$$

Here, $h_{t}: \mathbb{R}^{n} \times \mathbb{R}^{q} \rightarrow \mathbb{R}, t \in T$, are given functions, $v_{t}, t \in T$, are uncertain parameters which belongs to compact sets $\mathcal{V}_{t} \subseteq \mathbb{R}^{q}$.

As we know, robust optimization [14-16] is an useful approach to solve optimization problems with uncertainty. Following robust optimization methodology, we usually associate UMP with its robust counterpart

$$
\begin{cases}\operatorname{Min}_{\mathbb{R}_{+}^{p}} & \left(f_{1}(x), \ldots, f_{p}(x)\right) \\ \text { s.t. } & h_{t}\left(x, v_{t}\right) \leq 0, \forall v_{t} \in \mathcal{V}_{t}, t \in T \\ & x \in \mathbb{R}^{n}\end{cases}
$$

Recently, following robust optimization methodology, many interesting results devoted to (UMP) and its generalizations have been obtained from several different perspectives. By using scalarizing methods and robust optimization, Lee and Lee [17] establish necessary optimality theorems for robust weakly and properly efficient solutions of a multiobjective optimization problem with uncertainty. By virtue of a new concept of generalized convexity and robust type constraint qualification conditions, Chen et al. [18] give some optimality conditions and duality results for an uncertain nonconvex and nonsmooth multi-objective optimization problem. Guo and Yu [19] obtain optimality conditions for robust approximate quasi-weakly efficient solutions for uncertain multi-objective convex optimization problems. By combining robust optimization and scalarization technique, Sun et al. [20] give some new characterizations of Wolfe type robust approximate duality and saddle point theorems for a nonsmooth robust multi-objective optimization problem. Sun et al. [21] investigate optimality conditions for robust $\epsilon$-quasi efficient solutions of a class of uncertain semi-infinite multi-objective optimization under some tools of non-smooth analysis and a new modified scalarization technique. In addition, nonsmooth robust $\epsilon$ duality properties and $\epsilon$-quasi saddle point theorems are also established. New results on optimality and duality results for uncertain multiobjective polynomial optimization problems are given in [22]. By using tangential subdifferential and robust optimization, Liu et al. [23] obtained some characterizations of robust optimal solution sets for nonconvex uncertain semi-infinite optimization problems.

On the other hand, the fractional multi-objective optimization problem is an important subclass of multi-objective optimization problems. In the last decades, a wide variety of interesting works devoted to fractional multi-objective optimization problems and its generalizations have been given, see, for example, [24-33]. We observe that there are some papers devoted to the study of uncertain fractional multi-objective optimization problems under a robust optimization approach. In [34], the authors study approximate optimality conditions and Wolfe-type robust approximate duality of robust approximate weakly efficient solutions for uncertain fractional multi-objective optimization problems. Li et al. [35] establish optimality theorems and robust duality properties for minimax convex-concave fractional optimization problems with uncertainty. Antczak [36] establish a new parametric approach for robust approximate quasi-efficient solutions of robust fractional multi-objective optimization problems. Feng and Sun [37] obtain some new results for robust weakly $\epsilon$-efficient solutions for an uncertain fractional multi-objective semi-infinite optimization by employing conjugate analysis. Very recently, by employing robust limiting constraint qualification conditions and generalized convexity assumptions, Thuy and Su [38] consider optimality conditions and duality results for nonsmooth fractional multi-objective semi-infinite optimization problems with uncertain data.

In this paper, our main concern is to give new duality results of robust $\epsilon$-quasiefficient solutions for fractional multi-objective semi-infinite optimization problems (UFP, for brevity) with uncertainty appearing in the constraint functions. We first introduce the robust counterpart model (RFP, for brevity) for UFP. Then, with the help of a robust-type subdifferential constraint qualification, we present a necessary approximate optimality condition for robust $\epsilon$-quasi-efficient solutions for (UFP). Subsequently, we introduce a Mond-Weir-type robust approximate dual problem of (UFP) based on the obtained necessary optimality conditions. Then, we investigate robust weak, strong and converselike duality results between them under a new assumption of generalized convex-inclusion for Lipschitz functions.

This paper is organized as follows. In Section 2, we first recall some basic concepts in nonsmooth analysis and present approximate optimality results for robust $\epsilon$-quasi-efficient solutions of (UFP). In Section 3, we introduce a Mond-Weir-type robust approximate dual
problem for (UFP), and establish the robust $\epsilon$-quasi duality results between them. As a special case, we also deal with robust $\epsilon$-quasi duality results of the uncertain multi-objective optimization problem (UMP) and its robust approximate dual problem.

## 2. Mathematical Preliminaries

In this paper, let us recall some concepts and preliminary results [39,40]. Let $\mathbb{R}^{p}$ be the $p$-dimensional Euclidean space. We use the notation $\|\cdot\|$ for the Euclidean norm for $\mathbb{R}^{p}$. The nonnegative orthant of $\mathbb{R}^{p}$ is defined by $\mathbb{R}_{+}^{p}:=\left\{x=\left(x_{1}, \ldots, x_{n}\right) \mid x_{k} \geq 0, k=1, \ldots, n\right\}$. We always use the symbol $\langle\cdot, \cdot\rangle$ for the inner product in $\mathbb{R}^{p}$. The closed unit ball of $\mathbb{R}^{p}$ is denoted by $\mathbb{B}^{*}$. For a nonempty infinite index set $T$, the linear space $\mathbb{R}^{(T)}[41]$ is denoted by
$\mathbb{R}^{(T)}:=\left\{\gamma_{T}=\left(\gamma_{t}\right)_{t \in T} \mid \gamma_{t}=0\right.$ for all $t \in T$ except for finitely many $\left.\gamma_{t} \neq 0\right\}$.
Let $\mathbb{R}_{+}^{(T)}$ be the nonnegative cone of $\mathbb{R}^{(T)}$, i.e.,

$$
\mathbb{R}_{+}^{(T)}:=\left\{\gamma_{T} \in \mathbb{R}^{(T)} \mid \gamma_{t} \geq 0, \forall t \in T\right\}
$$

Let $\phi: \mathbb{R}^{p} \rightarrow \mathbb{R}$ be a locally Lipschitz function. The Clarke generalized directional derivative of $\phi$ at $x \in \mathbb{R}^{p}$ in the direction $d \in \mathbb{R}^{p}$ is defined by

$$
\phi^{c}(x ; d):=\limsup _{y \rightarrow x, t \downarrow 0} \frac{\phi(y+t d)-\phi(y)}{t} .
$$

The one-sided directional derivative of $\phi$ at $x \in \mathbb{R}^{p}$ in direction $d \in \mathbb{R}^{p}$ is defined by

$$
\phi^{\prime}(x ; d):=\lim _{t \downarrow 0} \frac{\phi(x+t d)-\phi(x)}{t} .
$$

We say that $\phi$ is quasidifferentiable at $x \in \mathbb{R}^{p}$ iff, for each $d \in \mathbb{R}^{n}, \varphi^{\prime}(x ; d)$ exists and $\varphi^{\prime}(x ; d)=\varphi^{c}(x ; d)$. The Clarke subdifferential $\partial^{c} \phi(x)$ of $\phi$ at $x \in \mathbb{R}^{p}$ is defined by

$$
\partial^{c} \phi(x):=\left\{\xi^{*} \in \mathbb{R}^{p} \mid \phi^{c}(x ; d) \geq\left\langle\xi^{*}, d\right\rangle, \forall d \in \mathbb{R}^{p}\right\}
$$

Obviously,

$$
\phi^{c}(x ; d)=\sup _{\xi \in \partial c}\langle\xi(x)<d\rangle, \quad \forall d \in \mathbb{R}^{n} .
$$

On the other hand, if $\phi: \mathbb{R}^{p} \rightarrow \mathbb{R}$ is a convex function, $\partial^{c} \phi(x)$ coincides with the convex subdifferential $\partial \phi(x)$, that is

$$
\partial \phi(x):=\left\{\tilde{\zeta}^{*} \in \mathbb{R}^{p} \mid \phi(y)-\phi(x) \geq\left\langle\xi^{*}, y-x\right\rangle, \forall y \in \mathbb{R}^{p}\right\}
$$

Let $\Omega \subseteq \mathbb{R}^{p}$ be a nonempty subset. The Clarke normal cone to $\Omega$ at $x \in \Omega$ is defined by

$$
N^{c}(\Omega, x):=\left\{\xi \in \mathbb{R}^{p} \mid\left\langle\zeta^{*}, w\right\rangle \leq 0, \forall w \in T_{\Omega}(x)\right\} .
$$

Here, $T_{\Omega}(x)$ is the Clarke tangent cone to $\Omega$ at $x \in \Omega$. Clearly, if $\Omega \subseteq \mathbb{R}^{n}$ is a nonempty closed convex set, $N^{c}(\Omega, x)$ becomes the following normal cone:

$$
N(\Omega, x):=\left\{\xi^{*} \in \mathbb{R}^{p} \mid\left\langle\xi^{*}, y-x\right\rangle \leq 0, \forall y \in \Omega\right\} .
$$

In what follows, let $f_{i}, g_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}, i=1, \ldots, p$, and $h_{t}: \mathbb{R}^{n} \rightarrow \mathbb{R}, t \in T$. We consider the following fractional multi-objective optimization problem

$$
\left\{\begin{array}{c}
\operatorname{Min}_{\mathbb{R}_{+}^{p}}\left(\frac{f_{1}(x)}{g_{1}(x)}, \ldots, \frac{f_{p}(x)}{g_{p}(x)}\right)  \tag{FP}\\
\text { s.t. } h_{t}(x) \leq 0, \forall t \in T, \\
\\
x \in \mathbb{R}^{n} .
\end{array}\right.
$$

The fractional optimization problem (FP) under uncertain data in the constraint functions becomes

$$
\begin{cases}\operatorname{Min}_{\mathbb{R}_{+}^{p}}\left(\frac{f_{1}(x)}{g_{1}(x)}, \ldots, \frac{f_{p}(x)}{g p(x)}\right)  \tag{UFP}\\ \text { s.t. } & h_{t}\left(x, v_{t}\right) \leq 0, \forall t \in T, \\ & x \in \mathbb{R}^{n} .\end{cases}
$$

Here $h_{t}: \mathbb{R}^{n} \times \mathbb{R}^{q} \rightarrow \mathbb{R} . v_{t} \in \mathcal{V}_{t} \subseteq \mathbb{R}^{q}, t \in T$ are uncertain parameters.
For (UFP), we consider its robust counterpart, namely
(RFP)

$$
\begin{cases}\operatorname{Min}_{\mathbb{R}_{+}^{p}}\left(\frac{f_{1}(x)}{g_{1}(x)}, \ldots, \frac{f_{p}(x)}{g_{p}(x)}\right) \\ \text { s.t. } & h_{t}\left(x, v_{t}\right) \leq 0, \forall v_{t} \in \mathcal{V}_{t}, t \in T, \\ & x \in \mathbb{R}^{n} .\end{cases}
$$

In this paper, without special statements, let $f_{i}, i=1, \ldots, p$, be locally Lipschitz functions with $f_{i}(x) \geq 0, \forall x \in \mathbb{R}^{n}$, and $g_{i}, i=1, \ldots, p$, be locally Lipschitz functions with $g_{i}(x)>0, \forall x \in \mathbb{R}^{n}$.

Now, we give the following important notations, which will be used later in this paper.
Definition 1. For (UFP). We say that $\mathcal{F}$ is the robust feasible set of (UFP) iff

$$
\mathcal{F}:=\left\{x \in \mathbb{R}^{n} \mid h_{t}\left(x, v_{t}\right) \leq 0, \forall v_{t} \in \mathcal{V}_{t}, t \in T\right\} .
$$

Now, we consider the concept of robust $\epsilon$-quasi efficient solution for (UFP). We refer the readers to $[19,21,37]$ for other kinds of robust approximate efficient solutions.

Definition 2. Let $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\} . \bar{x} \in \mathcal{F}$ is a robust $\epsilon$-quasi efficient solution of (UFP) if there is not $x \in \mathcal{F}$, such that

$$
\frac{f_{i}(x)}{g_{i}(x)} \leq \frac{f_{i}(\bar{x})}{g_{i}(\bar{x})}-\epsilon_{i}\|x-\bar{x}\|, \text { for all } i=1, \ldots, p
$$

and

$$
\frac{f_{j}(x)}{g_{j}(x)}<\frac{f_{j}(\bar{x})}{g_{j}(\bar{x})}-\epsilon_{j}\|x-\bar{x}\|, \text { for some } j \in\{1, \ldots, p\}
$$

Remark 1. Note that $g_{i} \equiv 1$, the concept of robust $\varepsilon$-quasi efficient solution of (UFP) deduces to the robust $\varepsilon$-quasi efficient solution of (UMP), i.e., there is not $x \in \mathcal{F}$, such that

$$
f_{i}(x) \leq f_{i}(\bar{x})-\epsilon_{i}\|x-\bar{x}\|, \text { for all } i=1, \ldots, p
$$

and

$$
f_{j}(x)<f_{j}(\bar{x})-\epsilon_{j}\|x-\bar{x}\|, \text { for some } j \in\{1, \ldots, p\} .
$$

For more details, see [20,21,42].
Definition 3 ([43] (Definition 3.2)). Consider (UFP). We say that the robust-type subdifferential constraint qualification condition RSCQ holds at $\bar{x} \in \mathcal{F}$, iff

$$
N^{c}(\mathcal{F}, \bar{x}) \subseteq \bigcup_{\substack{T_{T} \in T(\bar{x}),, v_{T} \in \mathcal{V}_{T}}}\left[\sum_{t \in T} \lambda_{t} \partial_{x}^{c} h_{t}\left(\bar{x}, v_{t}\right)\right],
$$

where $T(\bar{x})=\left\{\lambda_{T} \in \mathbb{R}_{+}^{(T)} \mid \lambda_{t} h_{t}\left(\bar{x}, v_{t}\right)=0, \forall v_{t} \in \mathcal{V}_{t}, t \in T\right\}$.
Next, we recall the following necessary optimality conditions for robust $\epsilon$-quasiefficient solutions for (UFP) under the RSCQ. For convenience, let $\epsilon:=\left(\epsilon_{1}, \ldots, \epsilon_{p}\right) \in$ $\mathbb{R}_{+}^{p} \backslash\{0\}$.

Proposition 1 ([44] (Theorem 1)). Let $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$. Assume that (RSCQ) holds at $\bar{x} \in \mathcal{F}$. If $\bar{x}$ is a robust $\epsilon$-quasi-efficient solution of (UFP), then there exist $\bar{\eta}_{t} \geq 0$, and $\bar{v}_{t} \in V_{t}, t \in T$, such that

$$
\begin{equation*}
0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(\bar{x})+\sum_{i=1}^{p} \phi_{i}(\bar{x}) \partial^{c}\left(-g_{i}\right)(\bar{x})+\sum_{t \in T} \bar{\eta}_{t} \partial^{c} h_{t}\left(\cdot, \bar{v}_{t}\right)(\bar{x})+2 \sum_{i=1}^{p} \varepsilon_{i} g_{i}(\bar{x}) \mathbb{B}^{*}, \tag{1}
\end{equation*}
$$

and

$$
\begin{equation*}
\bar{\eta}_{t} h_{t}\left(\bar{x}, \bar{v}_{t}\right)=0, \forall t \in T \tag{2}
\end{equation*}
$$

Here, $\phi_{i}(\cdot)=\frac{f_{i}(\bar{x})}{g_{i}(\bar{x})}-\epsilon_{i}\|\cdot-\bar{x}\|, i=1, \ldots, p$.
Remark 2. Proposition 1 extends [45] (Theorem 3.1) from the case of scalar optimization to the multi-objective setting.

In the case that $g_{i} \equiv 1$, the following result can be easily obtained by Proposition 1 .
Proposition 2. Let $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$. Assume that (RSCQ) holds at $\bar{x} \in \mathcal{F}$. If $\bar{x}$ is a robust $\epsilon$-quasiefficient solution of (UMP), then there exist $\bar{\eta}_{t} \geq 0$, and $\bar{v}_{t} \in V_{t}, t \in T$, such that

$$
\begin{equation*}
0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(\bar{x})+\sum_{t \in T} \bar{\eta}_{t} \partial^{c} h_{t}\left(\cdot, \bar{v}_{t}\right)(\bar{x})+2 \sum_{i=1}^{p} \varepsilon_{i} \mathbb{B}^{*}, \tag{3}
\end{equation*}
$$

and

$$
\begin{equation*}
\bar{\eta}_{t} h_{t}\left(\bar{x}, \bar{v}_{t}\right)=0, \forall t \in T \tag{4}
\end{equation*}
$$

## 3. Main Results

In this section, based on the optimality conditions obtained in Proposition 1, we establish a robust Mond-Weir-type approximate dual problem for (UMFP), and then investigate robust duality properties between them. Here, we only consider their robust $\epsilon$-quasiefficient solutions. For the sake of convenience in the sequel, we set $f:=\left(f_{1}, \ldots, f_{p}\right), g:=$ $\left(g_{1}, \ldots, g_{p}\right), h_{T}:=\left(h_{t}\right)_{t \in T}, \eta_{T}:=\left(\eta_{t}\right)_{t} \in \mathbb{R}_{+}^{(T)}, \mathcal{V}_{T}:=\prod_{t \in T} \mathcal{V}_{t}$, and $v_{T}:=\left(v_{t}\right)_{t \in T} \in \mathcal{V}_{T}$.

Let $y \in \mathbb{R}^{n}$ and $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$. For given $v_{t} \in \mathcal{V}_{t}, t \in T$, the Mond-Weir-type uncertain approximate dual problem (UFD) of (UFP) is
(UFD)

$$
\begin{cases}\operatorname{Max}_{\mathbb{R}_{+}^{p}}\left(\frac{f_{1}(y)}{g_{1}(y)}, \ldots, \frac{f_{p}(y)}{g_{p}(y)}\right) \\ \text { s.t. } & 0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(y)+\sum_{i=1}^{p} \frac{f_{i}(y)}{g_{i}(y)} \partial^{c}\left(-g_{i}\right)(y)+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)+2 \sum_{i=1}^{p} \epsilon_{i} g_{i}(y) \mathbb{B}^{*}, \\ & \eta_{t} h_{t}\left(y, v_{t}\right) \geq 0, t \in T, \\ & y \in \mathbb{R}^{n}, \epsilon_{i} \geq 0, i=1, \ldots, p, \eta_{t} \geq 0, t \in T .\end{cases}
$$

The optimistic counterpart of (UFD) is defined by


Here, the maximization is also over all the parameters $v_{t} \in \mathcal{V}_{t}, t \in T$. The feasible set of (OFD) is defined as

$$
\begin{gathered}
\widehat{\mathcal{F}}:=\left\{\left(y, \eta_{T}, v_{T}\right) \in \mathbb{R}^{n} \times \mathbb{R}_{+}^{(T)} \times \mathcal{V}_{T} \left\lvert\, 0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(y)+\sum_{i=1}^{p} \frac{f_{i}(y)}{g_{i}(y)} \partial^{c}\left(-g_{i}\right)(y)+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)\right.\right. \\
\left.+2 \sum_{i=1}^{p} \epsilon_{i} g_{i}(y) \mathbb{B}^{*}, \eta_{t} h_{t}\left(y, v_{t}\right) \geq 0, t \in T\right\}
\end{gathered}
$$

Remark 3. (i) Obviously, if $g_{i}(x) \equiv 1, i=1, \ldots, p$, (UFD) becomes the following conventional Mond-Weir-type uncertain approximate dual problem of (UMP)

$$
(\mathrm{UMD}) \quad \begin{cases}\operatorname{Max}_{\mathbb{R}_{+}^{p}}\left(f_{1}(y), \ldots, f_{p}(y)\right) \\ \text { s.t. } & 0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(y)+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)+2 \sum_{i=1}^{p} \epsilon_{i} \mathbb{B}^{*}, \\ & \eta_{t} h_{t}\left(y, v_{t}\right) \geq 0, t \in T \\ & y \in \mathbb{R}^{n}, \epsilon_{i} \geq 0, i=1, \ldots, p, \eta_{t} \geq 0, t \in T\end{cases}
$$

and (OFD) becomes the following Mond-Weir-type optimistic dual problem of (UMP)

$$
(\mathrm{OMD}) \quad \begin{cases}\operatorname{Max}_{\mathbb{R}_{+}^{p}}\left(f_{1}(y), \ldots, f_{p}(y)\right) \\ \text { s.t. } & 0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(y)+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)+2 \sum_{i=1}^{p} \epsilon_{i} \mathbb{B}^{*}, \\ & \eta_{t} h_{t}\left(y, v_{t}\right) \geq 0, t \in T \\ & y \in \mathbb{R}^{n}, \epsilon_{i} \geq 0, i=1, \ldots, p, \eta_{t} \geq 0, v_{t} \in \mathcal{V}_{t}, t \in T\end{cases}
$$

Here, we denote the feasible set of (OMD) by

$$
\begin{aligned}
\overline{\mathcal{F}}:=\left\{\left(y, \eta_{T}, v_{T}\right)\right. & \in \mathbb{R}^{n} \times \mathbb{R}_{+}^{(T)} \times \mathcal{V}_{T} \mid 0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(y)+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y) \\
& \left.+2 \sum_{i=1}^{p} \epsilon_{i} g_{i}(y) \mathbb{B}^{*}, \eta_{t} h_{t}\left(y, v_{t}\right) \geq 0, t \in T\right\}
\end{aligned}
$$

(ii) In the case that $\epsilon=0$ and there is no uncertainty in the constraint functions. Then, (UFP) becomes (FP), and (OMD) collapses to

$$
\begin{cases}\operatorname{Max}_{\mathbb{R}_{+}^{p}}\left(\frac{f_{1}(y)}{g_{1}(y)}, \ldots, \frac{f_{p}(y)}{g_{p}(y)}\right) \\ \text { s.t. } & 0 \in \sum_{i=1}^{p} \partial f_{i}(y)+\sum_{i=1}^{p} \frac{f_{i}(y)}{g_{i}(y)} \partial^{c}\left(-g_{i}\right)(y)+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}(y) \\ & \eta_{t} h_{t}(y) \geq 0, t \in T \\ & y \in \mathbb{R}^{n}, \eta_{t} \geq 0, t \in T\end{cases}
$$

Now, similar to Definition 2, we introduce robust $\epsilon$-quasi efficient solutions for (UFD).
Definition 4. Let $\varepsilon \in \mathbb{R}_{+}^{p} \backslash\{0\} .\left(\bar{y}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \widehat{\mathcal{F}}$ is said to be a robust $\varepsilon$-quasi efficient solution of (UFD), iff it is an $\varepsilon$-quasi efficient solution of (OFD), i.e., there is no $\left(y, \eta_{T}, v_{T}\right) \in \widehat{\mathcal{F}}$, such that

$$
\frac{f_{i}(y)}{g_{i}(y)} \geq \frac{f_{i}(\bar{y})}{g_{i}(\bar{y})}+\epsilon_{i}\|y-\bar{y}\|, \text { for all } i=1, \ldots, p
$$

and

$$
\frac{f_{j}(y)}{g_{j}(y)}>\frac{f_{j}(\bar{y})}{g_{j}(\bar{y})}+\epsilon_{j}\|y-\bar{y}\|, \text { for some } j \in\{1, \ldots, p\}
$$

Remark 4. In particular, if $g_{i} \equiv 1$, the concept of robust $\varepsilon$-quasi efficient solution of (UFD) deduces to the robust $\varepsilon$-quasi efficient solution of (UMD), i.e., there is no $\left(y, \eta_{T}, v_{T}\right) \in \overline{\mathcal{F}}$, such that

$$
f_{i}(y) \geq f_{i}(\bar{y})+\epsilon_{i}\|y-\bar{y}\|, \text { for all } i=1, \ldots, p
$$

and

$$
f_{j}(y)>f_{j}(\bar{y})+\epsilon_{j}\|y-\bar{y}\|, \text { for some } j \in\{1, \ldots, p\} .
$$

In order to give robust duality relations for (UFP) and (UFD), we introduce the new definition of generalized convex-inclusion for Lipschitz functions, which is inspired by [32] (Definition 3.4) and [21] (Definition 3.3).

Definition 5. Let $\Omega \subseteq \mathbb{R}^{n}$. $\left(f,-g, h_{T}\right)$ is said to generalized convex-inclusion on $\Omega$ at $x \in \Omega$, iff for any $y \in \Omega, \xi_{i}^{*} \in \partial^{c} f_{i}(x), \xi_{i}^{* *} \in \partial^{c}\left(-g_{i}\right)(x), i=1, \ldots, p$, and $\gamma_{t}^{*} \in \partial_{x}^{c} h_{t}\left(x, v_{t}\right), v_{t} \in \mathcal{V}_{t}$, $t \in T$, there exists $\omega \in \mathbb{R}^{n}$, such that

$$
\begin{gathered}
f_{i}(y)-f_{i}(x)>\left\langle\xi_{i}^{*}, \omega\right\rangle, i=1, \ldots, p, \\
-g_{i}(y)+g_{i}(x) \geq\left\langle\xi_{i}^{* *}, \omega\right\rangle, i=1, \ldots, p \\
h_{t}\left(y, v_{t}\right)-h_{t}\left(x, v_{t}\right) \geq\left\langle\gamma_{t}^{*}, \omega\right\rangle, t \in T, \\
\left\langle b^{*}, \omega\right\rangle \leq\|y-x\|, \forall b^{*} \in \mathbb{B}^{*},
\end{gathered}
$$

and

$$
0 \in \partial^{c} g_{i}(y), i=1, \ldots, p
$$

Remark 5. (i) In the special case that $g_{i} \equiv 1$, the concept of generalized convex-inclusion reduces to the concept of generalized convexity, i.e., $\left(f, h_{T}\right)$ is generalized convex on $\Omega$ at $x \in \Omega$, iff for any $y \in \Omega, \xi_{i}^{*} \in \partial^{c} f_{i}(x), i=1, \ldots, p$, and $\gamma_{t}^{*} \in \partial_{x}^{c} g_{t}\left(x, v_{t}\right), v_{t} \in \mathcal{V}_{t}, t \in T$, there exists $\omega \in \mathbb{R}^{n}$, such that

$$
\begin{aligned}
& f_{i}(y)-f_{i}(x)>\left\langle\xi_{i}^{*}, \omega\right\rangle, i=1, \ldots, p \\
& h_{t}\left(y, v_{t}\right)-h_{t}\left(x, v_{t}\right) \geq\left\langle\gamma_{t}^{*}, \omega\right\rangle, t \in T,
\end{aligned}
$$

and

$$
\left\langle b^{*}, \omega\right\rangle \leq\|y-x\|, \forall b^{*} \in \mathbb{B}^{*}
$$

(ii) If $g_{i} \equiv 1$ and there is uncertain data on $f_{i}, i=1, \ldots, p$, Definition 5 reduces to [21] (Definition 3.3).
(iii) If $g_{i} \equiv 1$ and there is no uncertain data on $h_{t}, t \in T$, Definition 5 reduces to the concept of generalized convexity-inclusion introduced in [32] (Definition 3.4), i.e., for any $y \in \Omega$, $\xi_{i}^{*} \in \partial^{c} f_{i}(x), \xi_{i}^{* *} \in \partial^{c}\left(-g_{i}\right)(x), i=1, \ldots, p$, and $\gamma_{t}^{*} \in \partial^{c} h_{t}(x), t \in T$, there exists $\omega \in \mathbb{R}^{n}$, such that

$$
\begin{gathered}
f_{i}(y)-f_{i}(x)>\left\langle\xi_{i}^{*}, \omega\right\rangle, i=1, \ldots, p, \\
-g_{i}(y)+g_{i}(x) \geq\left\langle\xi_{i}^{*}, \omega\right\rangle, i=1, \ldots, p, \\
h_{t}(y)-h_{t}(x) \geq\left\langle\gamma_{t}^{*}, \omega\right\rangle, t \in T, \\
\left\langle b^{*}, \omega\right\rangle \leq\|y-x\|, \forall b^{*} \in \mathbb{B}^{*},
\end{gathered}
$$

and

$$
0 \in \partial^{c} g_{i}(y), i=1, \ldots, p
$$

Note that this concept has been used to establish sufficient optimality conditions for weakly $\epsilon$-quasi-efficient solution for fractional optimization problem. For more details, please see [32] (Theorem 3.5).

Now, we show robust approximate duality properties for (UFP) and (UFD) by showing approximate duality properties between the robust counterpart (RMP) and the optimistic counterpart (OFD). In what follows, we set

$$
\begin{aligned}
& \omega_{1} \preceq \omega_{2} \Leftrightarrow \omega_{2}-\omega_{1} \in \mathbb{R}_{+}^{p} \backslash\{0\}, \forall \omega_{1}, \omega_{2} \in \mathbb{R}^{p} \\
& \omega_{1} \npreceq \omega_{2} \Leftrightarrow \omega_{2}-\omega_{1} \notin \mathbb{R}_{+}^{p} \backslash\{0\}, \forall \omega_{1}, \omega_{2} \in \mathbb{R}^{p} .
\end{aligned}
$$

The following result gives robust $\epsilon$-quasi-weak duality between (UFP) and (UFD).
Theorem 1. Let $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$. Suppose that $x \in \mathcal{F}$ and $\left(y, \eta_{T}, v_{T}\right) \in \widehat{\mathcal{F}}$. If $\left(f,-g, h_{T}\right)$ is generalized convex-inclusion on $\mathbb{R}^{n}$ at $y \in \mathbb{R}^{n}$, then,

$$
\left(\frac{f_{1}(x)}{g_{1}(x)}, \ldots, \frac{f_{p}(x)}{g_{p}(x)}\right) \npreceq\left(\frac{f_{1}(y)}{g_{1}(y)}-2 \epsilon_{1}\|x-y\|, \ldots, \frac{f_{p}(y)}{g_{p}(y)}-2 \epsilon_{p}\|x-y\|\right) .
$$

Proof. Suppose to the contrary that

$$
\left(\frac{f_{1}(x)}{g_{1}(x)}, \ldots, \frac{f_{p}(x)}{g_{p}(x)}\right) \preceq\left(\frac{f_{1}(y)}{g_{1}(y)}-2 \epsilon_{1}\|y-x\|, \ldots, \frac{f_{p}(y)}{g_{p}(y)}-2 \epsilon_{p}\|y-x\|\right) .
$$

Then,

$$
\begin{equation*}
\frac{f_{i}(x)}{g_{i}(x)} \leq \frac{f_{i}(y)}{g_{i}(y)}-2 \epsilon_{i}\|y-x\|, \text { for all } i=1, \ldots, p \tag{5}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{f_{i}(x)}{g_{i}(x)}<\frac{f_{j}(y)}{g_{j}(y)}-2 \epsilon_{j}\|y-x\|, \text { for some } j \in\{1, \ldots, p\} \tag{6}
\end{equation*}
$$

On the other hand, note that $\left(y, \eta_{T}, v_{T}\right) \in \widehat{\mathcal{F}}$. Then, $y \in \mathbb{R}^{n}, \eta_{t} \geq 0, v_{t} \in \mathcal{V}_{t}, t \in T$, and

$$
\begin{equation*}
0 \in \sum_{i=1}^{p} \partial^{c} f_{i}(y)+\sum_{i=1}^{p} \frac{f_{i}(y)}{g_{i}(y)} \partial^{c}\left(-g_{i}\right)(y)+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)+2 \sum_{i=1}^{p} \epsilon_{i} g_{i}(y) \mathbb{B}^{*}, \tag{7}
\end{equation*}
$$

and

$$
\begin{equation*}
\eta_{t} h_{t}\left(y, v_{t}\right) \geq 0, t \in T \tag{8}
\end{equation*}
$$

By (5), there exist $\xi_{i}^{*} \in \partial^{c} f_{i}(y), \zeta_{i}^{* *} \in \partial^{c}\left(-g_{i}\right)(y), i=1, \ldots, p, \zeta_{t}^{*} \in \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)$, $t \in T$, and $b^{*} \in \mathbb{B}^{*}$, such that

$$
\begin{equation*}
\sum_{i=1}^{p} \xi_{i}^{*}+\sum_{i=1}^{p} \frac{f_{i}(y)}{g_{i}(y)} \xi_{i}^{* *}+\sum_{t \in T} \eta_{t} \zeta_{t}^{*}+2 \sum_{i=1}^{p} \epsilon_{i} g_{i}(y) b^{*}=0 \tag{9}
\end{equation*}
$$

Since $\left(f,-g, h_{T}\right)$ is generalized convex-inclusion on $\mathbb{R}^{n}$ at $y \in \mathbb{R}^{n}$, we have for such $\xi_{i}^{*} \in \partial^{c} f_{i}(y), \xi_{i}^{* *} \in \partial^{c}\left(-g_{i}\right)(y), i=1, \ldots, p$, and $\zeta_{t}^{*} \in \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y), t \in T$, there exists $\vartheta \in \mathbb{R}^{n}$, such that

$$
f_{i}(x)-f_{i}(y)>\left\langle\xi_{i}^{*}, \vartheta\right\rangle, \quad i=1, \ldots, p
$$

$$
\begin{gathered}
-g_{i}(x)+g_{i}(y) \geq\left\langle\xi_{i}^{* *}, \vartheta\right\rangle, \quad i=1, \ldots, p \\
h_{t}\left(x, v_{t}\right)-h_{t}\left(y, v_{t}\right) \geq\left\langle\zeta_{t}^{*}, \vartheta\right\rangle, \quad t \in T, \\
\left\langle b^{*}, \vartheta\right\rangle \leq\|x-y\|, \forall b^{*} \in \mathbb{B}^{*},
\end{gathered}
$$

and

$$
0 \in \partial^{c} g_{i}(y), i=1, \ldots, p
$$

Together with (7)-(9), these follow that

$$
\begin{aligned}
& \sum_{i=1}^{p}\left(f_{i}(x)-\frac{f_{i}(y)}{g_{i}(y)} g_{i}(x)+2 \epsilon_{i} g_{i}(y)\|y-x\|\right) \\
& >\sum_{i=1}^{p}\left(f_{i}(y)+\left\langle\xi_{i}^{*}, \vartheta\right\rangle-\frac{f_{i}(y)}{g_{i}(y)} g_{i}(y)+\frac{f_{i}(y)}{g_{i}(y)}\left\langle\xi_{i}^{* *}, \vartheta\right\rangle+2 \epsilon_{i} g_{i}(y)\left\langle b^{*}, \vartheta\right\rangle\right) \\
& =\left\langle\sum_{i=1}^{p} \xi_{i}^{*}+\sum_{i=1}^{p} \frac{f_{i}(y)}{g_{i}(y)} \xi_{i}^{* *}+2 \sum_{i=1}^{p} \epsilon_{i} g_{i}(y) b^{*}, \vartheta\right\rangle \\
& =-\left\langle\sum_{i=1}^{p} \eta_{t} \zeta_{t}^{*}, \vartheta\right\rangle \\
& \geq-\sum_{t \in T} \eta_{t} h_{t}\left(x, v_{t}\right)+\sum_{t \in T} \eta_{t} h_{t}\left(y, v_{t}\right)
\end{aligned}
$$

Together with $\eta_{t} h_{t}\left(x, v_{t}\right) \leq 0, \forall x \in F$, and $\eta_{t} h_{t}\left(y, v_{t} \geq 0\right.$, we have

$$
\sum_{i=1}^{p}\left(f_{i}(x)-\frac{f_{i}(y)}{g_{i}(y)} g_{i}(x)+2 \epsilon_{i} g_{i}(y)\|y-x\|\right)>0
$$

Then, there exists $i_{0} \in\{1, \ldots, p\}$, such that

$$
f_{i_{0}}(x)-\frac{f_{i_{0}}(y)}{g_{i_{0}}(y)} g_{i_{0}}(x)+2 \epsilon_{i_{0}} g_{i_{0}}(y)\|y-x\|>0
$$

which follows that

$$
\begin{equation*}
\frac{f_{i_{0}}(x)}{g_{i_{0}}(x)}-\frac{f_{i_{0}}(y)}{g_{i_{0}}(y)}+2 \epsilon_{i_{0}} \frac{g_{i_{0}}(y)}{g_{i_{0}}(x)}\|y-x\|>0 \tag{10}
\end{equation*}
$$

Moreover, it follows from $0 \in \partial^{c} g_{i}(y), i=1, \ldots, p$, that

$$
\begin{equation*}
g_{i_{0}}(x) \geq g_{i_{0}}(y) \tag{11}
\end{equation*}
$$

Together with (10) and (11) , we have

$$
\frac{f_{i_{0}}(x)}{g_{i_{0}}(x)}-\frac{f_{i_{0}}(y)}{g_{i_{0}}(y)}+2 \epsilon_{i_{0}}\|y-x\|>0
$$

which is a contradiction to (5) and (6). Thus, the conclusion holds.
Now, we give the following example to justify the importance of the assumption of generalized convex-inclusion in Theorem 1.

Example 1. Let $\mathcal{V}_{t}:=[1-t, 1+t], t \in T:=\left[0, \frac{1}{2}\right]$. Let $f_{1}, f_{2}, g_{1}, g_{2}: \mathbb{R} \rightarrow \mathbb{R}$ and $g_{t}:$ $\mathbb{R} \times \mathbb{R} \rightarrow \mathbb{R}, t \in T$, be defined by

$$
f_{1}(x)=f_{2}(x):=\frac{1}{2}|x|+\frac{1}{6} x^{3}, g_{1}(x)=g_{2}(x):=|x|+1
$$

and

$$
h_{t}\left(x, v_{t}\right):=t x^{2}-t x-2 v_{t},
$$

where $x \in \mathbb{R}$ and $v_{t} \in \mathcal{V}_{t}, t \in T$. Then, (UFP) becomes

$$
\left\{\begin{array}{l}
\operatorname{Min}_{\mathbb{R}_{+}^{2}}\left(\frac{\frac{1}{2}|x|+\frac{1}{6} x^{3}}{|x|+1}, \frac{1}{2}|x|+\frac{1}{6} x^{3}\right. \\
|x|+1
\end{array}\right), ~ \begin{array}{ll}
\text { s.t. } \quad t x^{2}-t x-2 v_{t} \leq 0, \forall t \in\left[0, \frac{1}{2}\right], \\
\quad x \in \mathbb{R},
\end{array}
$$

and (RFP) becomes

$$
\left\{\begin{array}{l}
\operatorname{Min}_{\mathbb{R}_{+}^{2}}\left(\frac{\frac{1}{2}|x|+\frac{1}{6} x^{3}}{|x|+1}, \frac{\frac{1}{2}|x|+\frac{1}{6} x^{3}}{|x|+1}\right) \\
\text { s.t. } \quad t x^{2}-t x-2 v_{t} \leq 0, \forall v_{t} \in[1-t, 1+t], t \in\left[0, \frac{1}{2}\right] \\
\quad x \in \mathbb{R} .
\end{array}\right.
$$

Obviously, $\mathcal{F}=[-1,2]$. Let us consider $\bar{x}:=-1 \in \mathcal{F}$. Then,

$$
\left(\frac{f_{1}(\bar{x})}{g_{1}(\bar{x})}, \frac{f_{2}(\bar{x})}{g_{2}(\bar{x})}\right)=\left(\frac{1}{6}, \frac{1}{6}\right) .
$$

Now, consider the dual problem (UFD). In this setting, (OFD) becomes

$$
\left\{\begin{aligned}
& \operatorname{Max}_{\mathbb{R}_{+}^{2}}\left(\frac{f_{1}(y)}{g_{1}(y)}, \frac{f_{2}(y)}{g_{2}(y)}\right) \\
& \text { s.t. } 0 \in \partial^{c} f_{1}(y)+\partial^{c} f_{2}(y)+\frac{f_{1}(y)}{g_{1}(y)} \partial^{c}\left(-g_{1}\right)(y)+\frac{f_{2}(y)}{g_{2}(y)} \partial^{c}\left(-g_{2}\right)(y) \\
&+\sum_{t \in T} \eta_{t} \partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)+2 \epsilon_{1} g_{1}(y) \mathbb{B}^{*}+2 \epsilon_{2} g_{2}(y) \mathbb{B}^{*}, \\
& \eta_{t} h_{t}\left(y, v_{t}\right) \geq 0, t \in\left[0, \frac{1}{2}\right], \\
& y \in \mathbb{R}, \epsilon_{1} \geq 0, \epsilon_{2} \geq 0, \eta_{t} \geq 0, v_{t} \in[1-t, 1+t], t \in\left[0, \frac{1}{2}\right] .
\end{aligned}\right.
$$

Clearly, for any $y \in \mathbb{R}$ and $v_{T} \in \mathcal{V}_{T}$, we have

$$
\begin{gathered}
\partial^{c} f_{1}(y)=\partial^{c} f_{2}(y)=\left[\frac{1}{2} y^{2}-\frac{1}{2}, \frac{1}{2} y^{2}+\frac{1}{2}\right], \\
\partial^{c}\left(-g_{1}\right)(y)=\partial^{c}\left(-g_{2}\right)(y)=[-1,1],
\end{gathered}
$$

and

$$
\partial^{c} h_{t}\left(\cdot, v_{t}\right)(y)=\{2 t y-t\}, \forall t \in T .
$$

By selecting $\bar{y}:=1, \bar{\eta}_{t}:=0$, and $\bar{v}_{t}:=-t$, we have

$$
\begin{aligned}
& \partial^{c} f_{1}(\bar{y})+\partial^{c} f_{2}(\bar{y})+\frac{f_{1}(\bar{y})}{g_{1}(\bar{y})} \partial^{c}\left(-g_{1}\right)(\bar{y})+\frac{f_{2}(\bar{y})}{g_{2}(\bar{y})} \partial^{c}\left(-g_{2}\right)(\bar{y}) \\
& \quad+\sum_{t \in T} \bar{\eta}_{t} \partial^{c} h_{t}\left(\cdot, \bar{v}_{t}\right)(\bar{y})+2 \epsilon_{1} g_{1}(\bar{y}) \mathbb{B}^{*}+2 \epsilon_{2} g_{2}(\bar{y}) \mathbb{B}^{*} \\
& = \\
& {\left[-4 \epsilon_{1}-4 \epsilon_{2}-\frac{1}{3}, 4 \epsilon_{1}+4 \epsilon_{2}+\frac{7}{3}\right],}
\end{aligned}
$$

and

$$
\bar{\eta}_{t} h_{t}\left(\bar{y}, \bar{v}_{t}\right) \geq 0, t \in\left[0, \frac{1}{2}\right] .
$$

These mean that $\left(\bar{y}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \widehat{\mathcal{F}}$.

Now, take an arbitrarily $\epsilon=\left(\epsilon_{1}, \epsilon_{2}\right) \in \mathbb{R}_{+}^{2} \backslash\{0\}$ such that $\varepsilon_{i}<\frac{1}{12}, i=1,2$. Clearly,

$$
\begin{aligned}
\left(\frac{f_{1}(\bar{y})}{g_{1}(\bar{y})}-2 \epsilon_{1}\|\bar{x}-\bar{y}\|, \frac{f_{2}(\bar{y})}{g_{2}(\bar{y})}-2 \epsilon_{2}\|\bar{x}-\bar{y}\|\right) & =\left(\frac{1}{3}-2 \varepsilon_{1}, \frac{1}{3}-2 \varepsilon_{2}\right) \\
& \succ\left(\frac{1}{6}, \frac{1}{6}\right)=\left(\frac{f_{1}(\bar{x})}{g_{1}(\bar{x})}, \frac{f_{2}(\bar{x})}{g_{2}(\bar{x})}\right) .
\end{aligned}
$$

Thus, Theorem 1 is not applicable since $\left(f,-g, h_{T}\right)$ is not generalized convex-inclusion at $\bar{y}$. To do this, by choosing $\bar{\xi}_{i}:=0 \in \partial^{c} f_{i}(\bar{y}), i=1,2$, we have

$$
f_{i}(\bar{x})-f_{i}(\bar{y})=-\frac{2}{3}<0=\left\langle\bar{\xi}_{k}, \omega\right\rangle, \forall \omega \in \mathbb{R}
$$

Similarly, we obtain the following robust weak duality between (UMP) and (UMD).
Corollary 1. Let $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$. Suppose that $x \in \mathcal{F}$ and $\left(y, \eta_{T}, v_{T}\right) \in \overline{\mathcal{F}}$. If $\left(f, h_{T}\right)$ is generalized convex on $\mathbb{R}^{n}$ at $y \in \mathbb{R}^{n}$, then,

$$
\left(f_{1}(x), \ldots, f_{p}(x)\right) \npreceq\left(f_{1}(y)-2 \epsilon_{1}\|x-y\|, \ldots, f_{p}(y)-2 \epsilon_{p}\|x-y\|\right) .
$$

Remark 6. Clearly, by virtue of Example 1, we can also illustrate that the assumption of generalized convexity imposed in Corollary 1 is indispensable.

Now, we give robust strong duality results between (UFP) and (UFD).
Theorem 2. Let $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$. Assume that (RSCQ) holds at $\bar{x} \in \mathcal{F}$. Suppose that $\left(f,-g, h_{T}\right)$ is generalized convex-inclusion on $\mathbb{R}^{n}$ at $y \in \mathbb{R}^{n}$. If $\bar{x}$ is a robust $\epsilon$-quasi-efficient solution of (UFP), then there exist $\bar{\eta}_{T} \in \mathbb{R}_{+}^{(T)}$ and $\bar{v}_{T} \in \mathcal{V}_{T}$, such that $\left(\bar{x}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \widehat{\mathcal{F}}$ is a robust $2 \epsilon$-quasi-efficient solution of (UFD).

Proof. Assume that $\bar{x} \in \mathcal{F}$ is a robust $\epsilon$-quasi-efficient solution of (UFP). By Theorem 1, there exist $\bar{\eta}_{t} \geq 0$, and $\bar{v}_{t} \in \mathcal{V}_{t}, t \in T$, such that

$$
\begin{equation*}
0 \in \sum_{i=1}^{p} \partial f_{i}(\bar{x})-\sum_{i=1}^{p} \phi_{i}(\bar{x}) \partial g_{i}(\bar{x})+\sum_{t \in T} \bar{\eta}_{t} \partial h_{t}\left(\cdot, \bar{v}_{t}\right)(\bar{x})+2 \sum_{i=1}^{p} \varepsilon_{i} g_{i}(\bar{x}) \mathbb{B}^{*} \tag{12}
\end{equation*}
$$

and

$$
\begin{equation*}
\bar{\eta}_{t} h_{t}\left(\bar{x}, \bar{v}_{t}\right)=0, \forall t \in T \tag{13}
\end{equation*}
$$

From (12), (13) and $\phi_{i}(\bar{x})=\frac{f_{i}(\bar{x})}{g_{i}(\bar{x})}$, we have

$$
\left(\bar{x}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \widehat{\mathcal{F}} .
$$

By Theorem 1 , for all $\left(y, \eta_{T}, v_{T}\right) \in \widehat{\mathcal{F}}$, we have

$$
\left(\frac{f_{1}(\bar{x})}{g_{1}(\bar{x})}, \ldots, \frac{f_{p}(\bar{x})}{g_{p}(\bar{x})}\right) \npreceq\left(\frac{f_{1}(y)}{g_{1}(y)}-2 \epsilon_{1}\|\bar{x}-y\|, \ldots, \frac{f_{p}(y)}{g_{p}(y)}-2 \epsilon_{p}\|\bar{x}-y\|\right) .
$$

Thus, $\left(\bar{x}, \bar{\eta}_{T}, \bar{v}_{T}\right)$ is a robust $2 \epsilon$-quasi-efficient solutions of (UFD). Thus, the conclusion holds.

Remark 7. In [32] (Theorem 4.2), the authors established duality properties for $\epsilon$-quasi-weakly efficient solutions between (FP) and its Mond Weir-type dual problem. Therefore, Theorem 2 encompasses [32] (Theorem 4.2), where the corresponding results were given in terms of the similar methods.

Similarly, we give robust strong duality properties for robust $\epsilon$-quasi efficient solutions between (UMP) and (UMD).

Corollary 2. Let $\epsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$. Assume that (RSCQ) holds at $\bar{x} \in \mathcal{F}$. Suppose that $\left(f, h_{T}\right)$ is generalized convex on $\mathbb{R}^{n}$ at $y \in \mathbb{R}^{n}$. If $\bar{x}$ is a robust $\epsilon$-quasi-efficient solution of (UMP), then there exist $\bar{\eta}_{T} \in \mathbb{R}_{+}^{(T)}$ and $\bar{v}_{T} \in \mathcal{V}_{T}$, such that $\left(\bar{x}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \overline{\mathcal{F}}$ is a robust $2 \epsilon$-quasi-efficient solution of (UMD).

Now, we give a robust converse-like duality property between (UFP) and (UFD).
Theorem 3. Let $\varepsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$ and $\left(\bar{x}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \widehat{\mathcal{F}}$. If $\left(f,-g, h_{T}\right)$ is generalized convex-inclusion on $\mathbb{R}^{n}$ at $\bar{x} \in \mathcal{F}$, then, $\bar{x} \in \mathcal{F}$ is a robust $2 \varepsilon$-quasi efficient solution of (UMP).

Proof. Sine $\left(\bar{x}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \widehat{\mathcal{F}}$ and $\left(f,-g, h_{T}\right)$ is generalized convex-inclusion on $\mathbb{R}^{n}$ at $\bar{x}$, it follows from Theorem 1 that

$$
\left(\frac{f_{1}(x)}{g_{1}(x)}, \ldots, \frac{f_{p}(x)}{g_{p}(x)}\right) \npreceq\left(\frac{f_{1}(\bar{x})}{g_{1}(\bar{x})}-2 \epsilon_{1}\|x-\bar{x}\|, \ldots, \left.\frac{f_{p}(\bar{x})}{g_{p}(\bar{x})}-2 \epsilon_{p} \| x-\bar{x} \right\rvert\,\right), \forall x \in \mathcal{F} .
$$

Therefore, $\bar{x} \in \mathcal{F}$ is a robust $2 \varepsilon$-quasi efficient solution of (UFP) and the proof is complete.

Remark 8. Note that the converse-like duality result obtained in Theorem 3 extends [32] (Theorem 4.4) from the deterministic (i.e., with singleton uncertainty sets) to the robust setting. Moreover, Theorem 3 extends [43] (Theorem 4.3) from the scalar case to the multi-objective setting.

Similarly, we have the following results for (UMP) and (UMD), which has been considered in [21] (Theorem 4.3).

Corollary 3. Let $\varepsilon \in \mathbb{R}_{+}^{p} \backslash\{0\}$ and $\left(\bar{x}, \bar{\eta}_{T}, \bar{v}_{T}\right) \in \overline{\mathcal{F}}$. If $\left(f, h_{T}\right)$ is generalized convex on $\mathbb{R}^{n}$ at $\bar{x} \in \mathcal{F}$, then, $\bar{x} \in \mathcal{F}$ is a robust $\varepsilon$-quasi efficient solution of (UMP).

## 4. Conclusions

In this paper, we consider robust $\varepsilon$-quasi-efficient solutions for a class of uncertain fractional optimization problems. By employing robust optimization and the obtained optimality conditions, a Mond-Weir-type robust dual problem for the fractional optimization problem is established. Then, we give robust $\varepsilon$-quasi-weak, strong and converse duality properties between them in terms of generalized convex-inclusion assumptions. We also show that the obtained results extend the corresponding results obtained in [21,32,37].

In the future, similar to [21,43], it is of interest to formulate Mixed-type robust approximate dual problem of uncertain fractional optimization problems, and study robust $\varepsilon$-quasi-weak, strong, and converse duality properties between them.

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