INTERNATIONAL JOURNAL OF

## MATHEMATICAL COMBINATORICS



EDITED BY

THE MADIS OF CHINESE ACADEMY OF SCIENCES

June, 2011

# International Journal of <br> <br> Mathematical Combinatorics 

 <br> <br> Mathematical Combinatorics}

Edited By

The Madis of Chinese Academy of Sciences

June, 2011

Aims and Scope: The International J.Mathematical Combinatorics (ISSN 1937-1055) is a fully refereed international journal, sponsored by the MADIS of Chinese Academy of Sciences and published in USA quarterly comprising 100-150 pages approx. per volume, which publishes original research papers and survey articles in all aspects of Smarandache multi-spaces, Smarandache geometries, mathematical combinatorics, non-euclidean geometry and topology and their applications to other sciences. Topics in detail to be covered are:

Smarandache multi-spaces with applications to other sciences, such as those of algebraic multi-systems, multi-metric spaces, $\cdot$, etc.. Smarandache geometries;

Differential Geometry; Geometry on manifolds;
Topological graphs; Algebraic graphs; Random graphs; Combinatorial maps; Graph and map enumeration; Combinatorial designs; Combinatorial enumeration;

Low Dimensional Topology; Differential Topology; Topology of Manifolds;
Geometrical aspects of Mathematical Physics and Relations with Manifold Topology;
Applications of Smarandache multi-spaces to theoretical physics; Applications of Combinatorics to mathematics and theoretical physics;

Mathematical theory on gravitational fields; Mathematical theory on parallel universes;
Other applications of Smarandache multi-space and combinatorics.
Generally, papers on mathematics with its applications not including in above topics are also welcome.

It is also available from the below international databases:
Serials Group/Editorial Department of EBSCO Publishing
10 Estes St. Ipswich, MA 01938-2106, USA
Tel.: (978) 356-6500, Ext. 2262 Fax: (978) 356-9371
http://www.ebsco.com/home/printsubs/priceproj.asp
and
Gale Directory of Publications and Broadcast Media, Gale, a part of Cengage Learning
27500 Drake Rd. Farmington Hills, MI 48331-3535, USA
Tel.: (248) 699-4253, ext. 1326; 1-800-347-GALE Fax: (248) 699-8075
http://www.gale.com
Indexing and Reviews: Mathematical Reviews(USA), Zentralblatt fur Mathematik(Germany), Referativnyi Zhurnal (Russia), Mathematika (Russia), Computing Review (USA), Institute for Scientific Information (PA, USA), Library of Congress Subject Headings (USA).

Subscription A subscription can be ordered by a mail or an email directly to

## Linfan Mao

The Editor-in-Chief of International Journal of Mathematical Combinatorics
Chinese Academy of Mathematics and System Science
Beijing, 100190, P.R.China
Email: maolinfan@163.com

Price: US $\$ 48.00$

## Editorial Board (2nd)

## Editor-in-Chief <br> Linfan MAO

Chinese Academy of Mathematics and System Science, P.R.China
and
Beijing University of Civil Engineering and Architecture, P.R.China

Email: maolinfan@163.com

## Editors

## S.Bhattacharya

Deakin University
Geelong Campus at Waurn Ponds
Australia
Email: Sukanto.Bhattacharya@Deakin.edu.au

## Junliang Cai

Beijing Normal University, P.R.China
Email: caijunliang@bnu.edu.cn
Yanxun Chang
Beijing Jiaotong University, P.R.China
Email: yxchang@center.njtu.edu.cn

## Shaofei Du

Capital Normal University, P.R.China
Email: dushf@mail.cnu.edu.cn

## Xiaodong Hu

Chinese Academy of Mathematics and System
Science, P.R.China
Email: xdhu@amss.ac.cn

## Yuanqiu Huang

Hunan Normal University, P.R.China
Email: hyqq@public.cs.hn.cn

## H.Iseri

Mansfield University, USA
Email: hiseri@mnsfld.edu

Xueliang Li
Nankai University, P.R.China
Email: lxl@nankai.edu.cn
Ion Patrascu
Fratii Buzesti National College
Craiova Romania
Han Ren
East China Normal University, P.R.China
Email: hren@math.ecnu.edu.cn
Tudor Sireteanu, Dinu Bratosin and Luige
Vladareanu
Institute of Solid Mechanics of Romanian Academy
Bucharest, Romania.

## Guohua Song

Beijing University of Civil Engineering and Ar-
chitecture, P.R.China
songguohua@bucea.edu.cn
W.B.Vasantha Kandasamy

Indian Institute of Technology, India
Email: vasantha@iitm.ac.in

## Mingyao Xu

Peking University, P.R.China
Email: xumy@math.pku.edu.cn
Guiying Yan
Chinese Academy of Mathematics and System
Science, P.R.China
Email: yanguiying@yahoo.com
Y. Zhang

Department of Computer Science
Georgia State University, Atlanta, USA

Try not to become a man of success but rather try to become a man of value.

By A. Einstein, an American theoretical physicist.

# Duality Theorems of Multiobjective Generalized Disjunctive Fuzzy Nonlinear Fractional Programming 

E.E.Ammar<br>(Department of Mathematics, Faculty of Science, Tanta University, Egypt)<br>E-mail: amr.saed@ymail.com


#### Abstract

This paper is concerned with the study of duality conditions to convex-concave generalized multiobjective fuzzy nonlinear fractional disjunctive programming problems for which the decision set is the union of a family of convex sets. The Lagrangian function for such problems is defined and the Kuhn-Tucker Saddle and Stationary points are characterized. In addition, some important theorems related to the Kuhn-Tucker problem for saddle and stationary points are established. Moreover, a general dual problem is formulated together with weak; strong and converse duality theorems are proved.


Key Words: Generalized multiobjective fractional programming; Disjunctive programming; Convexity; Concavity; fuzzy parameters Duality.

AMS(2010): 49K45

## §1. Introduction

Fractional programming models have been became a subject of wide interest since they provide a universal apparatus for a wide class of models in corporate planning, agricultural planning, public policy decision making, and financial analysis of a firm, marine transportation, health care, educational planning, and bank balance sheet management. However, as is obvious, just considering one criterion at a time usually does not represent real life problems well because almost always two or more objectives are associated with a problem. Generally, objectives conflict with each other; therefore, one cannot optimize all objectives simultaneously. Nondifferentiable fractional programming problems play a very important role in formulating the set of most preferred solutions and a decision maker can select the optimal solution.

Chang in [8] gave an approximate approach for solving fractional programming with absolutevalue functions. Chen in [10] introduced higher-order symmetric duality in non-differentiable multiobjective programming problems. Benson in [6] studied two global optimization problems, each of which involves maximizing a ratio of two convex functions, where at least one of the two convex functions is quadratic form. Frenk in [12] gives some general results of the above Benson problem. The Karush-Kuhn-Tucker conditions in an optimization problem with interval-valued objective function are derived by Wu in [29].

[^0]Balas introduced Disjunctive programs in [3, 4,]. The convex hull of the feasible points has been characterized for these programs with a class of problems that subsumes pure mixed integer programs and for many other non-convex programming problems in [5]. Helbig presented in [17, 18] optimality criteria for disjunctive optimization problems with some of their applications. Gugat studied in $[15,16]$ an optimization a problem having convex objective functions, whose solution set is the union of a family of convex sets. Grossmann proposed in [14] a convex nonlinear relaxation of the nonlinear convex disjunctive programming problem. Some topics of optimizing disjunctive constraint functions were introduced in [28] by Sherali. In [7], Ceria studied the problem of finding the minimum of a convex function on the closure of the convex hull of the union of a finite number of closed convex sets. The dual of the disjunctive linear fractional programming problem was studied by Patkar in [25]. Eremin introduced in [11] disjunctive Lagrangian function and gave sufficient conditions for optimality in terms of their saddle points. A duality theory for disjunctive linear programming problems of a special type was suggested by Gon?alves in [13].

Liang In [21] gave sufficient optimality conditions for the generalized convex fractional programming. Yang introduced in [30] two dual models for a generalized fractional programming problem. Optimality conditions and duality were considered in [23] for nondifferentiable, multiobjective programming problems and in [20, 22] for nondifferentiable, nonlinear fractional programming problems. Jain et al in [19] studied the solution of a generalized fractional programming problem. Optimality conditions in generalized fractional programming involving nonsmooth Lipschitz functions are established by Liu in [23]. Roubi [26] proposed an algorithm to solve generalized fractional programming problem. Xu [31] presented two duality models for a generalized fractional programming and established its duality theorems. The necessary and sufficient optimality conditions to nonlinear fractional disjunctive programming problems for which the decision set is the union of a family of convex sets were introduced in [1]. Optimality conditions and duality for nonlinear fractional disjunctive minimax programming problems were considered in [2]. In this paper we define the Langrangian function for the nonlinear generalized disjunctive multiobjective fractional programming problem and investigate optimality conditions. For this class of problems, the Mond-Weir and Schaible type of duality are proposed. Weak, strong and converse duality theorems are established for each dual problem.

## §2. Problem Statement

Assume that $N=\{1,2, \cdots, p\}$ and $\mathcal{K}=\{1,2, \cdots, q\}$ are arbitrary nonempty index sets. For $i \in N$, let $g_{j}^{i}: \mathbf{R}^{n} \rightarrow \mathbf{R}$ be a vector map whose components are convex functions, $g_{j}^{i}(x) \leq 0$, $1 \leq j \leq m$. Suppose that $f_{r}^{i k}, h_{r}^{i+m+k}: \mathbf{R}^{n+q} \rightarrow \mathbf{r}$ are convex and concave functions for $i \in N, k \in \mathcal{K}, r=1, \cdots, s$ respectively, and $h_{r}^{i k}\left(x, \widetilde{b}_{r}\right)>0$. Here, these $\widetilde{a}_{r}, \widetilde{b}_{r}, r=1,2, \cdots, m$ represent the vectors of fuzzy parameters in the objectives functions. These fuzzy parameters are assumed to be characterized as fuzzy numbers [4].

We consider the generalized disjunctive multiobjective convex-concave fractional program
problem as in the following form:

$$
\begin{align*}
& \operatorname{GDFFVOP}(i) \quad \inf _{x \in Z_{i}} \max _{k \in K}\left\{\frac{f_{r}^{i k}\left(x, \widetilde{a}_{r}\right)}{h_{r}^{i k}\left(x, \widetilde{b}_{r}\right)}, r=1,2, \cdots, s\right\},  \tag{1}\\
& \text { Subject to } \quad \mathrm{x} \in \mathrm{Z}_{\mathrm{i}}, \mathrm{i} \in \mathrm{~N}, \tag{2}
\end{align*}
$$

where $Z_{i}=\left\{x \in \mathbf{R}^{n}: g_{j}^{i}(x) \leq 0, j=1,2, \cdots, m\right\}$. Assume that $Z_{i} \neq \emptyset$ for $i \in N$.
Definition $1([1])$ The $\alpha$-level set of the fuzzy numbers $\widetilde{a}$ and $\widetilde{b}$ are defined as the ordinary set $S_{\alpha}(\widetilde{a}, \widetilde{b})$ for which the degree of their membership functions exceeds level $\alpha$ :

$$
S_{\alpha}(\widetilde{a}, \widetilde{b})=\left\{(a, b) \in \mathbf{R}^{2 m} \mid \mu_{a r}\left(a_{r}\right) \geq \alpha, r=1,2, \cdots, m\right\} .
$$

For a certain degree of $\alpha$, the GDFVOP(i) problem can be written in the ordinary following form [11].

Lemma $1([7])$ Let $\alpha^{k}, \beta^{k}, k \in \mathcal{K}$ be real numbers and $\alpha^{k}>0$ for each $k \in K$. Then

$$
\begin{equation*}
\max _{k \in K} \frac{\beta^{k}}{\alpha^{k}} \geq \frac{\sum_{k \in K} \beta^{k}}{\sum_{k \in K} \alpha^{k}} \tag{3}
\end{equation*}
$$

By using Lemma 1 and from [9] The generalized multiobjective fuzzy fractional problem GDFFVOP(i) may be reformulated [3] as in the following two forms:
$\operatorname{GDFFNLP}(i, t, \alpha)$ :

$$
\begin{equation*}
\inf _{i \in N} \inf _{x \in Z_{i}(S)}\left\{\frac{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)},\left(a_{r}, b_{r}\right) \in S_{\alpha}(\widetilde{a}, \widetilde{b}), r=1,2, \cdots, m\right\} \tag{4}
\end{equation*}
$$

where $t^{k} \in \mathbf{R}_{+}^{q}$. Denote by

$$
M_{i}=\inf _{x \in Z_{i}} \frac{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)},\left(a_{r}, b_{r}\right) \in S_{\alpha}(\widetilde{a}, \widetilde{b}), r=1,2, \cdots, m
$$

the minimal value of $\operatorname{GDFFNLP}(i, t, \alpha)$, and let

$$
P_{i}=\left\{x \in Z_{i}: \frac{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)}=M_{i}, i \in N\right\}
$$

be the set of solutions of $\operatorname{GDFFNLP}(i, t, \alpha)$. The generalized multiobjective disjunctive fuzzy fractional programming problem is formulated as:

$$
\begin{equation*}
\operatorname{GDFFNLP}(t, \alpha): \quad \inf _{i \in N} \inf _{x \in Z}\left\{\frac{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)}\right\} \tag{5}
\end{equation*}
$$

where $t^{k} \in \mathbf{R}_{+}^{q}, k \in \mathcal{K}$ and $Z=\bigcup_{i \in N} Z_{i}$ is the feasible solution set of problem $\operatorname{GDFFNLP}(t, \alpha)$. For problem $\operatorname{GDFFNLP}(t, \alpha)$, we assume the following sets:
(I) $M=\inf _{i \in N} M_{i}$ is the minimal value of $\operatorname{GDFFNLP}(t, \alpha)$.
(II) $Z^{*}=\left\{x \in Z: \exists i \in I(X), \inf _{i} \frac{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)}=M\right\}$ is set of these of solutions on the problem $\operatorname{GDFFNLP}(t, \alpha)$, where $I=\left\{i \in I^{\prime}: x \in Z\right\}, I^{\prime}=\left\{i \in N: Z^{*} \neq \emptyset\right\}$ and $I^{\prime}=\{1,2, \cdots, a\} \subset N$. Problem $\operatorname{GDFFVOP}(t, \alpha)$ may be reformulated in the following form:
$\operatorname{GDFFNLP}(t, \alpha, d)$ :

$$
\begin{equation*}
\inf _{i \in I} \inf _{x \in Z}\left\{F^{i}\left(x, t, d^{i}, a, b\right)=\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)-d^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)\right\} \tag{6}
\end{equation*}
$$

where

$$
d^{i}=\frac{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)}>0, \quad i \in I
$$

We define the Lagrangian functions of problems $\operatorname{GDFFNLP}(t, \alpha, d)$ and $\operatorname{GDFFNLP}(t, \alpha)[21$, 24, and 25] in the following forms:

$$
\begin{equation*}
G L^{i}\left(x, \lambda^{i}, a, b\right)=F^{i}\left(x, t, d^{i}, a, b\right)+\lambda \sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(x) \tag{7}
\end{equation*}
$$

and

$$
\begin{equation*}
L^{i}\left(x, u, \lambda^{i}, a, b\right)=\frac{u^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)+\sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(x)}{u^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)} \tag{8}
\end{equation*}
$$

where $\lambda_{j}^{i} \geq 0$ and $u^{i} \geq 0, i \in I$ are Lagrangian multipliers. Then the Lagrangian functions $G L(x, \lambda, a, b)$ and $L(x, u, \lambda, a, b)$ of $\operatorname{GDFFNLP}(t, \alpha, d)$ are defined by:

$$
\begin{equation*}
G L(x, \lambda, a, b)=\inf _{i \in I} G L^{i}\left(x, \lambda^{i}, a, b\right)=\inf _{i \in I}\left\{F^{i}\left(x, t, d^{i}, a, b\right)+\sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(x)\right\} \tag{9}
\end{equation*}
$$

and

$$
\begin{equation*}
L(x, u, \lambda, a, b)=\inf _{i \in I} L^{i}\left(x, u, \lambda^{i}, a, b\right)=\inf _{i \in I}\left\{\frac{u^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}\left(x, a_{r}\right)+\sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(x)}{u^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}\left(x, b_{r}\right)}\right\} \tag{10}
\end{equation*}
$$

where $x \in Z, t^{k} \in \mathbf{R}_{+}^{q}, u \in \mathbf{R}_{+}^{q}$ and $\lambda \in \mathbf{R}_{+}^{q}$ are Lagrangian multipliers, respectively.

## §3. Optimality Theorems with Differentiability

Definition 3.1 A point $\left(x^{0}, \lambda^{0}, a^{)}, b^{)}\right)$in $R^{n+p+2 m}$ with $\lambda^{)} \geq 0$ is said to be a GL-saddle point of problem $\operatorname{GDFFNLP}(t, \alpha, d)$ if and only if

$$
\begin{equation*}
G L\left(x^{0}, \lambda, a, b\right) \leq G L\left(x^{0}, \lambda^{0}, a^{0}, b^{0}\right) \leq G L\left(x, \lambda^{0}, a^{0}, b^{0}\right) \tag{11}
\end{equation*}
$$

for all with $x \in \mathbf{R}^{n+p}$ and $\lambda \in \mathbf{R}_{+}^{m}$.
Definition 3.1 A point $\left({ }^{0}, u^{0}, \lambda^{0}\right)$ in $\mathbf{R}^{n+p+m}$, with $u^{0} \geq 0$ and $\lambda^{0} \geq 0$ is said to be an L-saddle point of problem $\operatorname{GDFFNLP}(t, \alpha)$ if and only if

$$
\begin{equation*}
L\left(x^{0}, \lambda, a, b\right) \leq L\left(x^{0}, \lambda^{0}, a^{0}, b^{0}\right) \leq L\left(x, \lambda^{0}, a^{0}, b^{0}\right) \tag{12}
\end{equation*}
$$

for all with $x \in \mathbf{R}^{n+p}, u \in \mathbf{R}_{+}^{m}$ and $\lambda \in \mathbf{R}_{+}^{m}$.
The proof of the following theorems follows as in [3].
Theorem 3.1(Sufficient Optimality Criteria) If for $d^{0 i} \geq 0$ the point $\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right)$ is a saddle point of $G L(x, \lambda, a, b)$ and $F^{i}\left(x, t, d^{0 i}, a^{0}, b^{0}\right), g_{j}^{i}(x)$ are bounded and convex functions.


Corollary 3.1 If the point $\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right)$ is a saddle point of $L(x, u, \mu)$ and $F^{i}\left(x, t, d^{i}, a^{0}, b^{0}\right)$, $g_{j}^{i}(x)$ are bounded and convex functions. Then $x^{0}$ is a minimal solution for the problem $G D F F N L P(t, \alpha)$.

The proof is follows similarly as proof of Theorem 3.1.
Assumption 3.1 Let $F^{i}\left(x, y, d^{i}, a, b\right)=0$ be a convex function on $\operatorname{Conv} Z\left(Z=\bigcup_{i \in I}\right)$. If for all $x \in \operatorname{Conv} Z$, the functions $F^{i}\left(x, t^{0}, d^{0 i}, a^{0}, b^{0}\right)-F^{i}\left(x^{0}, t^{0}, d^{0 i}, a^{0}, b\right), x^{0} \in \operatorname{Conv} Z, i \in I$, $t^{0} \in \mathbf{R}_{+}^{q}$ and $\left(a^{0}, b^{0}\right) \in \mathbf{R}^{2 m}$ are bounded, then $\inf _{i \in I}\left\{F^{i}\left(x, t^{0}, d^{0 i}, a^{0}, b\right)-F^{i}\left(x^{0}, t^{0}, d^{0 i}, a^{0}, b\right)\right\}$ is a convex function on Conv $Z$.

Proposition 3.1 Under the Assumption 3.1, and if the system

$$
\left.\begin{array}{l}
\inf _{i \in I} F^{i}\left(x, t^{0}, d^{0 i}, a^{0}, b\right)-F^{i}\left(x^{0}, t^{0}, d^{0 i}, a^{0}, b^{0}\right)<0 \\
g_{j}^{i}(x) \leq 0 \text { for at least one } i \in I
\end{array}\right\}
$$

has no solution on Conv $Z$, then $\exists \lambda^{0} \in \mathbf{R}_{+}, \lambda^{0 i} \in \mathbf{R}_{+}^{m},\left(\lambda^{0}, \lambda^{0 i}\right) \geq 0$ and $t^{0} \in \mathbf{R}_{+}^{q}$ such that

$$
\mu^{0} \inf _{i \in i} F^{i}\left(x, t^{0}, d^{0 i}, a^{0}, b^{0}\right)+\inf _{i \in i} \sum_{j=1}^{m} \mu_{j}^{0 i} g_{j}^{i}(x) \geq 0
$$

for $\forall x \in \operatorname{Conv} Z$.
Corollary 3.2 With Assumption 3.1, $g_{j}^{i}(x), i \in I, j=1,2, \cdots, m$ satisfy the $C Q$ and $x^{0}$ is an optimal solution of problem $\operatorname{GDFFNLP}(t, \alpha)$, then there exists $u^{0} \geq 0$ and $\lambda^{0} \geq 0$ such that $\left(x^{0}, t^{0}, \lambda^{0}, a^{0}, b^{0}\right)$ is a saddle point of $L\left(x^{0}, t^{0}, \lambda^{0}, a^{0}, b^{0}\right)$.

The proof is follows similarly as proof of Theorem 3.2.

## §4. Optimality Theorems without Differentiability

Definition 4.1 The point $\left(x^{0}, \lambda^{0}, a^{0}, b^{0}\right), x^{0} \in x \in \mathbf{R}^{n+p}, \lambda^{0}, a^{0}, b^{0} \in \mathbf{R}^{3 m}$, if they exist such that

$$
\begin{array}{lc}
\nabla_{x} G L\left(x^{0}, \lambda^{0}, a^{0}, b^{0}\right) \geq 0, & x^{0} \nabla_{x} G L\left(x^{0}, \lambda^{0}, a^{0}, b^{0}\right)=0 \\
\nabla_{\lambda x} G L\left(x^{0}, \lambda^{0}, a^{0}, b^{0}\right) \geq 0, & \lambda^{0} \nabla_{\lambda} G L\left(x^{0}, \lambda^{0}, a^{0}, b^{0}\right)=0 \\
\sum_{j=1}^{m} \lambda_{j}^{0 i} g_{j}^{i}\left(x^{0}\right)=0, \quad \lambda_{j}^{i} \geq 0, & i \in I, \quad j=1,2, \cdots, m \tag{15}
\end{array}
$$

is could Kuhn- Tucker stationary point of problem $\operatorname{GDFFNLP}\left(t^{0}, \alpha^{0}, d^{0}\right)$. Or, equivalently,

$$
\begin{align*}
& \nabla_{x} \inf _{i \in I}\left\{F^{i}\left(x^{0}, t^{0}, d^{0 i}, a^{0}, b^{0}\right)+\lambda_{j}^{0 i} g_{j}^{i}\left(x^{0}\right)\right\}=0, \quad i \in I  \tag{16}\\
& g_{j}^{i}\left(x^{0}\right) \leq 0, \quad i \in I, \quad j=1,2, \cdots, m  \tag{17}\\
& \sum_{j=1}^{m} \lambda_{j}^{0 i} g_{j}^{i}\left(x^{0}\right)=0, \quad \lambda_{j}^{i} \geq 0, \quad i \in I, \quad j=1,2, \cdots, m \tag{18}
\end{align*}
$$

Definition 4.2 The point $\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right), x \in \mathbf{R}^{n+p+2 m}, u \in \mathbf{R}_{+}^{q}$ and $\lambda \in \mathbf{R}_{+}^{m}$, if they exist such that

$$
\begin{array}{ll}
\nabla_{x} L\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right) \geq 0, & x^{0} \nabla_{x} L\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right)=0 \\
\nabla_{u} L\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right) \geq 0, & u^{0} \nabla_{\lambda} L\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right)=0 \\
\nabla_{\mu} L\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right) \geq 0, & \mu^{0} \nabla_{\lambda} L\left(x^{0}, u^{0}, \lambda^{0}, a^{0}, b^{0}\right)=0 \\
\sum_{j=1}^{m} \lambda_{j}^{0 i} g_{j}^{i}\left(x^{0}\right)=0, \quad \lambda_{j}^{i} \geq 0, & i \in I, \quad j=1,2, \cdots, m \tag{22}
\end{array}
$$

is could Kuhn- Tucker stationary point of problem $\operatorname{GDFFNLP}\left(t^{0}, \alpha^{0}\right)$. Or, equivalently,

$$
\begin{align*}
& \nabla_{x} \inf _{i \in I}\left\{\frac{u^{0 i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{0 k} f_{r}^{i k}\left(x^{0}, a^{0}\right)+\sum_{j=1}^{m} \lambda_{j}^{0 i} g_{j}^{i}\left(x^{0}\right)}{u^{0 i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{0 k} h_{r}^{i k}\left(x^{0}, b^{0}\right)}\right\}=0  \tag{23}\\
& g_{j}^{i}\left(x^{0}\right) \leq 0, \quad i \in I, \quad j=1,2, \cdots, m  \tag{24}\\
& \sum_{j=1}^{m} \lambda_{j}^{0 i} g_{j}^{i}\left(x^{0}\right)=0, \quad \lambda_{j}^{i} \geq 0, \quad i \in I, \quad j=1,2, \cdots, m \tag{25}
\end{align*}
$$

The proof of the following theorem follows as in [3].
Theorem 4.1 Assume that $F^{i}\left(x, t, d^{i}, a, b\right), g_{j}^{i}(x), i \in I, j=1,2, \cdots, m$ are convex differentiable functions on Conv $S$. If $F^{i}\left(x, t, d^{i}, a, b\right)$ and $g_{j}^{i}(x)$ are bounded functions for each $x \in C o v S$ and $g_{j}^{i}(x)$ satisfy $C Q$ for $i \in I$, then $x^{0}$ is an optimal solution of $G D F F N L P(t, \alpha, d)$ if and only if there are Lagrange multipliers $\lambda^{0} \in \mathbf{R}^{p+m}, \lambda \geq 0$ such that (13)-(15) are satisfied.

Corollary 4.1 Suppose that $F^{i}\left(x, t, d^{i}, a, b\right), g_{j}^{i}(x), i \in I, j=1,2, \cdots, m$ are convex differentiable functions on Conv S. If $F^{i}\left(x, t, d^{i}, a, b\right)$ and $g_{j}^{i}(x)$ are bounded functions for each $x \in \operatorname{Cov} S$ and $g_{j}^{i}(x)$ satisfy $C Q$ for $i \in I$, then $x^{0}$ is an optimal solution of $\operatorname{GDFFNLP}(t, \alpha, d)$ if and only if there are Lagrange multipliers $u^{0} \geq 0, \lambda \geq 0, u \in \mathbf{R}_{+}^{q}$ and $\lambda^{0} \in \mathbf{R}^{p+m}$ such that (19)-(22) are satisfied.

The proof is follows similarly as the proof of Theorem 4.1.

Theorem 4.2 Assume that $F^{i}\left(x, t, d^{i}, a, b\right)$ is a pseudoconvex function at $x \in C o n v S$ and that $\sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(x)$ is a quasiconvex function. If $F^{i}\left(x, t, d^{i}\right)$ and $g_{j}^{i}(x)$ are bounded functions for each $x \in \operatorname{Conv} S$, and if the equations (28)-(30) are satisfied for $\operatorname{tin} \mathbf{R}_{+}^{k}$ and $\lambda^{0} \in \mathbf{R}_{+}^{p+m}$, then $x^{0}$ is an optimal solution of $\operatorname{GDFFNL}(t, \alpha, d)$.

Corollary 4.2 Assume that $F^{i}\left(x, t, d^{i}, a, b\right)$ is a pseudoconvex function at $x \in C o n v S$ and that $\sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(x)$ is a quasiconvex function. If $F^{i}\left(x, t, d^{i}, a, b\right)$ and $g_{j}^{i}(x)$ are bounded functions for each $x \in \operatorname{Conv} S$ and there exists $u^{0} \in \mathbf{R}_{+}^{k}$ and $\lambda^{0} \in \mathbf{R}_{+}^{p+m}$ such that equations (16)-(18) are satisfied, then $x^{0}$ is an optimal solution of $\operatorname{GDFFNLP}(t, \alpha, d)$.

The proof is follows similarly as proof of Theorem 4.2.

## §5. Duality Using Mond-Weir Type

According to optimality Theorems 4.1 and 4.2 , we can formulate the Mond-Weir type dual (M-WDGF) of the disjunctive fractional minimax problem $\operatorname{GDFFNLP}(t, \alpha, d)$ as follows:
$\mathrm{M}-\mathrm{WDGF} \max _{y \in \mathbf{R}^{n}} \sup _{i \in I}\left(H^{i}(y, t, \alpha, D, a, b)=\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(y, a)-D^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(y, b)\right)$,
where

$$
D^{i}=\frac{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(y, a)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(y, b)}>0, \quad i \in I
$$

Problem (M-WDGF) satisfies the following conditions:

$$
\begin{align*}
& \sup _{i \in I} \nabla_{y}\left\{H(y, t, D, a, b)+\sum \lambda_{j}^{i} g_{j}^{i}(y)\right\}=0  \tag{27}\\
& \sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(y)=0, \quad \lambda_{j}^{i} \geq 0, \quad i \in I, \quad j=1,2, \cdots, m  \tag{28}\\
& \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(y, a)-D^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(y, b) \geq 0, \quad i \in I, \quad D^{i}>0 \tag{29}
\end{align*}
$$

Theorem 5.1(Weak Duality) Let $x$ be feasible for $\operatorname{GDFFNLP}(t, \alpha, d)$ and ( $u, \lambda, t, a, b)$ be feasible for (M-WDGFD). If for all feasible $(y, \lambda, t, a, b), H^{i}(y, t, \alpha, D, a, b)$ are pseudoconvex
for each $i \in I$, and $\sum_{j=1}^{m} \lambda_{j}^{i} g_{j}^{i}(y)$ are quasiconvex for $i \in I$, then $\inf (G D F F N L P(t, \alpha, d)) \geq$ $\sup (M-W D G F)$.

Proof If not, then there must be that

$$
\begin{aligned}
& \inf _{i \in I} \inf _{x \in Z}\left(\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(x)-d^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(x)\right) \\
& <\sup _{i \in I} \sup _{y \in \mathbf{R}^{n}}\left(\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(y)-D^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(y)\right) .
\end{aligned}
$$

Hence, for $i \in I$, we get that

$$
\begin{equation*}
\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(x)-d^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(x)<\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(y)-D^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(y) . \tag{30}
\end{equation*}
$$

and by the pseudoconvexity of $H^{i}(y, t, D),(30)$ implies that

$$
\begin{equation*}
(x-y)^{t} \nabla_{x}\left(\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(y)-D^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(y)\right)<0 \tag{31}
\end{equation*}
$$

Equation (31) implies that

$$
\begin{equation*}
\sup _{i \in I} \sup _{y \in \mathbf{R}^{n}}\left((x-y)^{t} \nabla_{x}\left(\sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} f_{r}^{i k}(y)-D^{i} \sum_{r=1}^{s} \sum_{k=1}^{K} t^{k} h_{r}^{i k}(y)\right)\right)<0 . \tag{32}
\end{equation*}
$$

From equation (27) and inequality (32) it follows that

$$
\begin{equation*}
\sup _{i \in I}\left\{(x-y)^{t} \nabla_{x} \sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(y)\right\}>0 . \tag{33}
\end{equation*}
$$

By (26), inequality (33) implies that

$$
\sup _{i \in I} \sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(x)>\sup _{i \in I} \sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(u)>0 .
$$

Then $\sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(x)>0$, which contradicts the assumption that $x$ is feasible with respect to $\operatorname{GDFFNLP}(t, \alpha, d)$.

Theorem 5.2(Strong Duality) If $x^{0}$ is an optimal solution of $\operatorname{GDFFNLP(t,\alpha ,d)\text {and}CQ}$ is satisfied, then there exists $\left(y^{0}, \lambda^{0}, t^{0}, a^{0}, b^{0}\right) \in \mathbf{R}^{n+m}$ is feasible for ( $M-W D G F$ ) and the corresponding value of $\inf (G D F F N L P(t, \alpha, d))=\sup (M-W D G F)$.

Proof Since $x^{0}$ is an optimal solution of $\operatorname{DGFFNLP}\left(t^{0}, \alpha^{0}, d^{0}\right)$ and satisfy CQ, then there is a positive integer $\lambda_{j}^{* i} \geq 0, i \in I, j=1,2, \cdots, m$ such that Kuhn-Tucker conditions (27)-(29)
are satisfied. Assume that $\lambda^{0}=\tau^{-1} \lambda^{*}$ in the Kuhn-Tucker stationary point conditions. It follows that $\left(y^{0}, \lambda^{0}, t^{0}, a^{0}, b^{0}\right)$ is feasible for (M-WDGF). Hence

$$
\inf _{i \in I}\left(\frac{\sum_{k=1}^{K} t^{0 k} f_{r}^{i k}\left(x^{0}, a^{0}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{0 k} h_{r}^{i k}\left(x^{0}, b^{0}\right)}\right)=\sup _{i \in I}\left(\frac{\sum_{k=1}^{K} t^{0 k} f_{r}^{i k}\left(y^{0}, a^{0}\right)}{\sum_{r=1}^{s} \sum_{k=1}^{K} t^{0 k} h_{r}^{i k}\left(y^{0}, b^{0}\right)}\right) .
$$

Theorem 5.3(Converse Duality) Let $x^{0}$ be an optimal solution of $\operatorname{DGFFNLP}\left(t^{0}, \alpha^{0}, d^{0}\right)$ and $C Q$ is satisfied. If $\left(y^{*}, \mu^{*}\right)$ is an optimal solution of ( $M-W D F D$ ) and $H^{i}\left(y^{*}, t^{*}, D^{*}\right)$ is strictly pseudoconvex at $y^{*}$, then $y^{*}=x^{0}$ is an optimal solution of $\operatorname{GDFFNLP}(t, \alpha, d)$.

Proof Let $x^{0}$ be an optimal solution of $\operatorname{DGFFNLP}\left(t^{0}, \alpha^{0}, d^{0}\right)$ and CQ is satisfied. Assume that $y^{*} \neq x^{0}$. Then $\left(y *, \mu^{*}\right)$ is an optimal solution of (M-WDGF). Whence,

$$
\begin{equation*}
\inf _{i \in I} \inf _{k \in K} F^{i}\left(x^{0}, t^{0}, d^{0 i}\right)=\sup _{i \in I} \sup _{k \in K} H^{i}\left(y^{*}, t^{*}, D^{* i}\right) \tag{34}
\end{equation*}
$$

Because $\left(y^{*}, \mu^{*}\right)$ is feasible with respect to (M-WDGF), it follows that

$$
\sum_{j=1}^{m} \mu_{j}^{* i} g_{j}^{i}\left(x^{0}\right) \leq \sum_{j=1}^{m} \mu_{j}^{* i} g_{j}^{i}\left(y^{*}\right)
$$

Quasiconvexity of $\sum_{j=1}^{m} \mu_{j}^{* i} g_{j}^{i}(x)$ implies that

$$
\begin{equation*}
\sup _{i \in i}\left(x^{0}-y^{*}\right) \sum_{j=1}^{m} \nabla_{x} \mu_{j}^{* i} g_{j}^{i}\left(y^{*}\right) \leq 0 \tag{35}
\end{equation*}
$$

From (34) and (35), it follows that

$$
\begin{equation*}
\sup _{i \in i}\left(x^{0}-y^{*}\right) \nabla_{y} H^{i}\left(y^{*}, t^{*}, D^{* i}\right) \geq 0 \tag{36}
\end{equation*}
$$

From (36) and the strict pseudoconvexity of at $y^{*}$, it follows that

$$
\sup _{i \in i} \nabla_{x} F^{i}\left(x^{0}, t^{0}, d^{0 i}\right)>\sup _{i \in i} \nabla_{y} H^{i}\left(y^{*}, t^{*}, D^{* i}\right) .
$$

This contradicts to (35). Hence $y^{*}=x^{0}$ is an optimal solution of $\operatorname{DGFFNLP}\left(t^{0}, \alpha^{0}, d^{0}\right)$.

## §6. Duality Using Schaible Formula

The Schaible dual of $\operatorname{GDFFNLP}(t, \alpha, d)$ has been formulated in [27] as follows:

$$
\begin{equation*}
\max _{(y, \mu) \in \mathbf{R}^{n+m}} D \tag{SGD}
\end{equation*}
$$

where $(y, \mu) \in \mathbf{R}^{n} \times \mathbf{R}_{+}^{m}$ satisfying:

$$
\begin{align*}
& \sup _{i \in I} \nabla_{x}\left\{\sum_{k=1}^{K} t^{k} f^{i k}(y)-D^{i} \sum_{k=1}^{K} t^{k} h^{i k}(y)+\sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(y)\right\}=0,  \tag{37}\\
& \sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(y) \geq 0, \quad i \in I,  \tag{38}\\
& \sum_{k=1}^{K} t^{k} f^{i k}(y)-D^{i} \sum_{k=1}^{K} t^{k} h^{i k}(y) \geq 0, \quad i \in i,  \tag{39}\\
& D^{i} \geq 0 \quad \text { and } \quad \mu_{j}^{i} \geq 0, \quad i \in I, \quad j=1,2, \cdots, m . \tag{40}
\end{align*}
$$

 feasible $(y, \mu)$, $\sup _{i \in I} H^{i}(y, t, d)$ is pseudoconvex at $u$ and $\sup _{i \in I} \sum_{j=1}^{K} \mu_{j}^{i} g_{j}^{i}(y)$ is quasiconvex, then $\inf G D F F N L P(t, \alpha, d) \geq \sup (S G D)$.

Proof For each $i \in I$, suppose that

$$
\frac{\sum_{k=1}^{K} t^{k} f^{i k}(y)}{\sum_{k=1}^{K} t^{k} h^{i k}(y)}<D^{i} .
$$

Hence, for each $y \in \mathbf{R}^{n}$ and $i \in I$, we get that

$$
\sum_{k=1}^{K} t^{k} f^{i k}(y)-D^{i} \sum_{k=1}^{K} t^{k} h^{i k}(y)<0
$$

Therefore,

$$
\begin{equation*}
\sup _{i \in I}\left(\sum_{k=1}^{K} t^{k} f^{i k}(y)-D^{i} \sum_{k=1}^{K} t^{k} h^{i k}(y)\right)<0 \tag{41}
\end{equation*}
$$

From (39) and (41) with $t \neq 0$, we have

$$
\left(\sum_{k=1}^{K} t^{k} f^{i k}(x)-d^{i} \sum_{k=1}^{K} t^{k} h^{i k}(x)\right)<\left(\sum_{k=1}^{K} t^{k} f^{i k}(y)-D^{i} \sum_{k=1}^{K} t^{k} h^{i k}(y)\right) .
$$

By the pseudoconvexity of $\sup _{i \in I} H^{i}(y, t, D)$ at $u$, it follows that

$$
\begin{equation*}
(x-y)^{T}\left(\sum_{k=1}^{K} t^{k} f^{i k}(y)-D^{i} \sum_{k=1}^{K} t^{k} h^{i k}(y)\right)<0 \tag{42}
\end{equation*}
$$

Consequently, (38) and (42) yield that

$$
\begin{equation*}
(x-y)^{T} \sum_{j=1}^{m} \mu_{j}^{i} \nabla_{x} g_{j}^{i}(y)>0 \tag{43}
\end{equation*}
$$

and, by the quasiconvexity of $\sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(y)$, inequality (43) implies that

$$
\begin{equation*}
\sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(x)>\sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(y) \tag{44}
\end{equation*}
$$

From inequalities (38) and (44) it follows that

$$
\begin{equation*}
\sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(x)>0 \tag{45}
\end{equation*}
$$

But, from the feasibility of $x \in S$ and $\mu_{j}^{i} \geq 0, i \in I, j=1,2, \cdots, m$, (1) implies that $\sum_{j=1}^{m} \mu_{j}^{i} g_{j}^{i}(x) \leq 0$, this contradicts (45). Hence,

$$
\frac{\sum_{k=1}^{K} t^{k} f^{i k}(y)}{\sum_{k=1}^{K} t^{k} h^{i k}(y)} \geq D^{i}
$$

i.e., $\inf G D F F N L P(t, \alpha, d) \geq \sup (S G D)$.

Theorem 6.2(Strong Duality) Let $x^{0}$ be an optimal solution of $\operatorname{GDFFNLP}(t, \alpha, d)$ so that $C Q$ is satisfied. Then there exists $\left(y^{0}, \mu^{0}\right)$ is feasible for $(S D D)$ and the corresponding value of $\inf G D F F N L P(t, \alpha, d)=\sup (S D D)$. If, in addition, the hypotheses of Theorem 6.1 are satisfied, then $\left(x^{0}, \mu^{0}\right)$ is an optimal solution of (SDD).

Proof The proof is similar to that of Theorem 5.2.
Theorem 6.3(Converse Duality) Suppose that $x$ ) is an optimal solution of GDFFNLP $(t, \alpha, d)$ and $g_{j}^{i}(x)$ satisfy CQ. Let the hypotheses of the above Theorem 6.1 hold. If $\left(y^{*}, \mu^{*}\right)$ is an optimal solution of (SDD) and is strictly pseudocovex at $y^{*}$, then $y^{*}=x^{0}$ is an optimal solution of $D G F F N L P\left(t^{0}, \alpha^{0}, d^{0}\right)$.

Proof Assume that $y^{*} \neq x^{0}, x^{0}$ is an optimal solution $\operatorname{DGFFNLP}\left(t^{0}, \alpha^{0}, d^{0}\right)$ and try to find a contraction. From Theorem 4.2, for each $i \in I$, it follows that

$$
\begin{equation*}
\frac{\sum_{k=1}^{K} t^{0 k} f^{i k}\left(x^{0}\right)}{\sum_{k=1}^{K} t^{0 k} h^{i k}\left(x^{0}\right)}=d^{0 i} \tag{46}
\end{equation*}
$$

Applying (1) with (38) we get that

$$
\sum_{j=1}^{m} \mu_{j}^{* i} g_{j}^{i}\left(x^{0}\right) \leq \sum_{j=1}^{m} \mu_{j}^{* i} g_{j}^{i}\left(y^{*}\right)
$$

By quasiconvexity of $\sum_{j=1}^{m} \mu_{j}^{* i} g_{j}^{i}(x)$ and for each $i \in I$, it follows that

$$
\begin{equation*}
\left(x^{0}-y^{*}\right) \sum_{j=1}^{m} \nabla_{x} \mu_{j}^{* i} g_{j}^{i}\left(y^{*}\right) \leq 0 \tag{47}
\end{equation*}
$$

From (37) and (47) it follows that

$$
\begin{equation*}
\left(x^{0}-y^{*}\right) \nabla_{x}\left(\sum_{k=1}^{K} t^{* k} f^{i k}\left(y^{*}\right)-D^{* i} \sum_{k=1}^{K} t^{* k} h^{i k}\left(y^{*}\right)\right) \leq 0 \tag{48}
\end{equation*}
$$

From (39), (48) and the strict pseudoconvexity of $\left(\sum_{k=1}^{K} t^{* k} f^{i k}(y)-D^{* i} \sum_{k=1}^{K} t^{* k} h^{i k}(y)\right)$ for each $i \in I$ at $y^{*}$, it follows that

$$
\begin{equation*}
\left(\sum_{k=1}^{K} t^{0 k} f^{i k}\left(x^{0}\right)-d^{0 i} \sum_{k=1}^{K} t^{0 k} h^{i k}\left(x^{0}\right)\right)>\left(\sum_{k=1}^{K} t^{* k} f^{i k}\left(y^{*}\right)-D^{* i} \sum_{k=1}^{K} t^{* k} h^{i k}\left(y^{*}\right)\right) \tag{49}
\end{equation*}
$$

Inequality (49) implies that

$$
\begin{equation*}
\left(\sum_{k=1}^{K} t^{0 k} f^{i k}(x)-d^{0 i} \sum_{k=1}^{K} t^{0 k} h^{i k}(x)\right)>0, \quad i \in I \tag{50}
\end{equation*}
$$

i.e., for each $i \in I$ it is follows that

$$
\begin{equation*}
\frac{\sum_{k=1}^{K} t^{0 k} f^{i k}(x)}{\sum_{k=1}^{K} t^{0 k} h^{i k}(x)}>d^{0 i} \tag{51}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
\frac{\sum_{k=1}^{K} t^{0 k} f^{i k}\left(x^{0}\right)}{\sum_{k=1}^{K} t^{0 k} h^{i k}\left(x^{0}\right)} \geq \frac{\sum_{k=1}^{K} t^{0 k} f^{i k}(x)}{\sum_{k=1}^{K} t^{0 k} h^{i k}(x)}>d^{0 i} \tag{52}
\end{equation*}
$$

contradicts to that (46). So that $y^{*}=x^{0}$ is an optimal solution of $\operatorname{DGFFNLP}\left(t^{0}, \alpha^{0}, d^{0}\right)$.

## $\S 7$. Conclusion

This paper addresses the solution of generalized multiobjective disjunctive programming problems, which corresponds to minmax continuous optimization problems that involve disjunctions with convex-concave nonlinear fractional objective functions. We use Dinkelbach's global approach for finding the maximum of this problem. We first describe the Kuhn-Tucker saddle point of nonlinear disjunctive fractional minmax programming problems by using the decision set that is the union of a family of convex sets. Also, we discuss necessary and sufficient optimality conditions for generalized nonlinear disjunctive fractional minmax programming problems. For the class of problems, we study two duals; we propose and prove weak, strong and converse duality theorems.

## References

[1] E. E.Ammar, On optimality and duality theorems of nonlinear disjunctive fractional minmax programs, Euro. J. Oper. Res. 180(2007) 971-982.
[2] E. E.Ammar, On optimality of nonlinear fractional disjunctive programming problem, Computer $\mathcal{E}$ Math. With Appl., USA, 53(2007), 1527-1537.
[3] E. E.Ammar, A study on optimality and duality theorems of nonlinear generalized disjunctive fractional programming, Math. and Computer Modeling, 48(2008), 69-82.
[3] E.Balas, Disjunctive programming, Annals of Discrete Mathematics, 5, 3-1, (1979).
[4] E.Balas, Disjunctive programming and a hierarchy of relaxation for discrete optimization problems, SIAMJ. Alg.Disc, Math., 6,466-486, (1985).
[5] E.Balas, Disjunctive programming: properties of the convex hull of feasible points, Disc. App., Math., 89, No. 1-3, 3-44 (1998).
[6] H.P.Benson, Fractional programming with convex quadratic forms and functions, Euro. J. Oper. Res., 173 (2)(2006) 351-369.
[7] S. Ceria and J.Soares, Convex programming for disjunctive convex optimization, Math. Program., 86A. No. 3, 595-614 (1999).
[8] Chang Ching-Ter, An approximate approach for fractional programming with absolutevalue functions, J. Appl. Math. Comput., 161, No. 1 (2005).171-179.
[9] V.Changkang and Y.Haimes, Multiobjective decision making, Theory and Methodology, Series, Vol.8, North-Holland, New York 1983.
[10] Chen Xiuhong, Higher-order symmetric duality in non-differentiable multiobjective programming problems, J. Math. Anal. Appl., 290, No.2, 423- 435 (2004).
[11] Eremin Ivan I., About the problem of disjunctive programming, J. Yugosl. J. Oper. Res., 10, No.2, 149-161 (2000).
[12] J.B.G.Frenk, A note on paper Fractional programming with convex quadratic forms and functions by Benson H.P., Euro. J. Oper. Res. 176(2007) 641-642.
[13] Goncalves, Amilcar S., Symmetric duality for disjunctive programming with absolute value functional, J. Eur. J. Oper. Res., 26, 301-306 (1986).
[14] I.E.Grossmann and S. Lee, Generalized convex disjunctive programming nonlinear convex hull relaxation, Comput. Optim. Appl., 26, No.1, 83-100(2003).
[15] M. Gugat, One - sided derivatives for the value function in convex parametric Programming, Optimization, Vol.28, 301-314(1994).
[16] M. Gugat, Convex parametric programming optimization One-sided differentiability of the value function, J. Optim. Theory and Appl. , Vol. 92, 285-310 (1997).
[17] S. Helbig, An algorithm for vector optimization problems over a disjunctive feasible set, Operation Research, 1989.
[18] S. Helbig, Duality in disjunctive programming via vector optimization, J. Math. Program, 65A, No.1, 21-41 (1994).
[19] S. Jain and A.Magal, Solution of a generalized fractional programming problem, J. Indian Acad. Math., 26, No.1, 15-21 (2004).
[20] D. S. Kim, S. J. Kim and M. H.Kim, Optimality and duality for a class of nondifferentiable multiobjective fractional programming problems, Journal of Optimization Theory
and Applications, 129 (2006) 131-146.
[21] Liang Zhi-an, Shi Zhen-wei, Optimality conditions and duality for a minimax fractional programming with generalized convexity, J. Math. Anal. Appl. 277, No.2, 474-488 (2003).
[22] S. Liu, E. Feng, Optimality conditions and duality for a class of non-differentiable nonlinear fractional programming problems, Int. J. Appl. Math., 13, No.4, 345-358 (2003).
[23] J.C. Liu, Y. Kimura and K. Tanaka, Generalized fractional programming, RIMS Kokyuroku, 1031, 1-13 (1998).
[24] S.K.Mishra, S.Y.Wang, K.K.Lai, Higher-order duality for a class of nondifferentiable multiobjective programming problems, Int. J. Pure Appl. Math., 11, No.2, 221-232 (2004).
[25] Patkar Vivek and I.M.Stancu-Minasian, Duality in disjunctive linear fractional programming, Eur. J. Oper. Res., 21, 101-105 (1985).
[26] A. Roubi, Method of centers for generalized fractional programming, J. Optim. Theory Appl., 107, No. 1 (2000) 123-143.
[27] S. Schaible, Fractional programming I, Duality, Mang. Sci., 12 (1976) 858-867.
[28] H.D. Sherali and C.M. Shetty, Optimization with Disjunctive Constraints, Springer Verlag Berlin, Heidelberg, New York, (1980).
[29] Wu H.-Ch., The Karush-Kuhn-Tucker optimality conditions in an optimization problem with interval-valued objective function, Euro. J. Oper. Res., 176(2007)46-59.
[30] X.M. Yang, X.Q.Yang and K.L.Teo, Duality and saddle- point type optimality for generalized nonlinear fractional programming, J. Math. Anal. Appl., 289, No.1, 100-109 (2004).
[31] Z.K. Xu, Duality in generalized nonlinear fractional programming, J. Math. Anal. Appl., 169(1992) 1-9.

# Surface Embeddability of Graphs via Joint Trees 

Yanpei Liu<br>(Institute of Mathematics, Beijing Jiaotong University, Beijing 100044, P.R.China)<br>E-mail: ypliu@bjtu.edu.cn


#### Abstract

This paper provides a way to observe embedings of a graph on surfaces based on join trees and then characterizations of orientable and nonorientable embeddabilities of a graph with given genus.


Key Words: Surface, graph, Smarandache $\lambda^{S}$-drawing, embedding, joint tree.
AMS(2010): 05C15, 05C25

## §1. Introduction

A drawing of a graph $G$ on a surface $S$ is such a drawing with no edge crosses itself, no adjacent edges cross each other, no two edges intersect more than once, and no three edges have a common point. A Smarandache $\lambda^{S}$-drawing of G on $S$ is a drawing of G on $S$ with minimal intersections $\lambda^{S}$. Particularly, a Smarandache 0-drawing of $G$ on $S$, if existing, is called an embedding of $G$ on $S$.

The term joint three looks firstly appeared in [1] and then in [2] in a certain detail and [3] firstly in English. However, the theoretical idea was initiated in early articles of the author [4-5] in which maximum genus of a graph in both orientable and nonorientable cases were investigated.

The central idea is to transform a problem related to embeddings of a graph on surfaces i.e., compact 2-manifolds without boundary in topology into that on polyhegons (or polygons of even size with binary boundaries). The following two principles can be seen in [3].

Principle A Joint trees of a graph have a 1-to-1 correspondence to embeddings of the graph with the same orientability and genus i.e., on the same surfaces.

Principle B Associate polyhegons (as surfaces) of a graph have a 1-to-1 correspondence to joint trees of the graph with the same orientability and genus, i.e., on the same surfaces.

The two principle above are employed in this paper as the theoretical foundation. These enable us to discuss in any way among associate polyhegons, joint trees and embeddings of a graph considered.

[^1]
## §2. Layers and Exchangers

Given a surface $S=(A)$. it is divided into segments layer by layer as in the following.
The 0 th layer contains only one segment, i.e., $A\left(=A_{0}\right)$;
The 1 st layer is obtained by dividing the segment $A_{0}$ into $l_{1}$ segments, i.e., $S=\left(A_{1}, A_{2}\right.$, $\cdots, A_{l_{1}}$ ), where $A_{1}, A_{2}, \cdots, A_{l_{1}}$ are called the 1 st layer segments;

Suppose that on $k-1$ st layer, the $k-1$ st layer segments are $A_{\underline{n}_{(k-1)}}$ where $\underline{n}_{(k-1)}$ is an integral $k-1$-vector satisfied by

$$
\underline{1}_{(k-1)} \leqslant\left(n_{1}, n_{2}, \cdots, n_{k-1}\right) \leqslant \underline{N}_{(k-1)}
$$

with $\underline{1}_{(k-1)}=(1,1, \cdots, 1), \underline{N}_{(k-1)}=\left(N_{1}, N_{2}, \cdots, N_{k-1}\right), N_{1}=l_{1}=N_{(1)}, N_{2}=l_{A_{N_{(1)}}}$, $N_{3}=l_{A_{\underline{N}_{(2)}}}, \cdots, N_{k-1}=l_{A_{\underline{N}_{(k-2)}}}$, then the $k$ th layer segments are obtained by dividing each $k-1$ st layer segment as

$$
\begin{equation*}
A_{\underline{\underline{n}}_{(k-1)}, 1}, A_{\underline{n}_{(k-1)}, 2}, \cdots, A_{\underline{\underline{n}}_{(k-1)},}, l_{\underline{\underline{n}}_{(k-1)}} \tag{1}
\end{equation*}
$$

where $\underline{1}_{(k)}=\left(\underline{n}_{(k-1)}, 1\right) \leqslant\left(\underline{n}_{(k-1)}, i\right) \leqslant \underline{N}_{(k)}=\left(\underline{N}_{(k-1)}, N_{k}\right)$ and $N_{k}=l_{\underline{A}_{(k-1)}}$. Segments in (1) are called successors of $A_{\underline{n}_{(k-1)}}$. Conversely, $A_{\underline{n}_{(k-1)}}$ is the predecessor of any one in (1).

A layer segment which has only one element is called an end segment and others, principle segments. For an example, let

$$
S=(1,-7,2,-5,3,-1,4,-6,5,-2,6,7,-3,-4)
$$

Fig.2.1 shows a layer division of $S$ and Tab.2.1, the principle segments in each layer.
For a layer division of a surface, if principle segments are dealt with vertices and edges are with the relationship between predecessor and successor, then what is obtained is a tree denoted by $T$. On $T$, by adding cotree edges as end segments, a graph $G=(V, E)$ is induced. For example, the graph induced from the layer division shown in Fig. 1 is as

$$
\begin{equation*}
V=\{A, B, C, D, E, F, G, H, I\} \tag{2}
\end{equation*}
$$

and

$$
\begin{equation*}
E=\{a, b, c, d, e, f, g, h, 1,2,3,4,5,6,7\}, \tag{3}
\end{equation*}
$$

where

$$
\begin{aligned}
& a=(A, B), b=(A, C), c=(A, D), d=(B, E), \\
& e=(C, F), f=(C, G), g=(D, H), h=(D, I),
\end{aligned}
$$

and

$$
\begin{array}{r}
1=(B, F), 2=(E, H), 3=(F, I), 4=(G, I), \\
5=(B, C), 6=(G, H), 7=(D, E) .
\end{array}
$$

By considering $E_{T}=\{a, b, c, d, e, f, g, h\}, \bar{E}_{T}=\{1,2,3,4,5,6,7\}, \delta_{i}=0, i=1,2, \cdots, 7$, and the rotation $\sigma$ implied in the layer division, a joint tree $\widehat{T}_{\sigma}^{\delta}$ is produced.


Fig. 1 Layer division of surface $S$

| Layers | Principle segments |
| ---: | ---: |
| 0th layer | $A=\langle 1,-7,2-5 ; 3,-1,4,-6,5 ;-2,6,7,-3-4\rangle$ |
| 1st layer | $B=\langle 1 ;-7,2 ;-5\rangle, C=\langle 3,-1 ; 4,-6 ; 5\rangle$, |
|  | $D=\langle-2,6 ; 7 ;-3,-4\rangle$ |
| 2nd layer | $E=\langle-7 ; 2\rangle, F=\langle 3 ;-1\rangle, G=\langle 4 ;-6\rangle$, |
|  | $H=\langle-2 ; 6\rangle, I=\langle-3 ;-4\rangle$ |

Tab. 1 Layers and principle segments
Theorem 1 A layer division of a polyhegon determines a joint tree. Conversely, a joint tree determines a layer division of its associate polyhegon.

Proof For a layer division of a polyhegon as a polyhegon, all segments are treated as vertices and two vertices have an edge if, and only if, they are in successive layers with one as a subsegment of the other. This graph can be shown as a tree. Because of each non-end vertex with a rotation and end vertices pairwise with binary indices, this tree itself is a joint tree.

Conversely, for a joint tree, it is also seen as a layer division of the surface determined by the boundary polyhegon of the tree.

Then, an operation on a layer division is discussed for transforming an associate polyhegon into another in order to visit all associate polyhegon without repetition.

A layer segment with all its successors is called a branch in the layer division. The operation of interchanging the positions of two layer segments with the same predecessor in a layer division is called an exchanger.

Lemma 1 A layer division of an associate polyhegon of a graph under an exchanger is still a layer division of another associate polyhegon. Conversely, the later under the same exchanger becomes the former.

Proof On the basis of Theorem 1, only necessary to see what happens by exchanger on a joint tree once. Because of only changing the rotation at a vertex for doing exchanger once,
exchanger transforms a joint tree into another joint tree of the same graph. This is the first conclusion. Because of exchanger inversible, the second conclusion holds.

On the basis of this lemma, an exchanger can be seen as an operation on the set of all associate surfaces of a graph.

Lemma 2 The exchanger is closed in the set of all associate polyhegons of a graph.
Proof From Theorem 1, the lemma is a direct conclusion of Lemma 1.
Lemma 3 Let $\mathcal{A}(G)$ be the set of all associate polyhegons of a graph $G$, then for any $S_{1}$, $S_{2} \in \mathcal{A}(G)$, there exist a sequence of exchangers on the set such that $S_{1}$ can be transformed into $S_{2}$.

Proof Because of exchanger corresponding to transposition of two elements in a rotation at a vertex, in virtue of permutation principle that any two rotation can be transformed from one into another by transpositions, from Theorem 1 and Lemma 1, the conclusion is done.

If $\mathcal{A}(G)$ is dealt as the vertex set and an edge as an exchanger, then what is obtained in this way is called the associate polyhegon graph of $G$, and denoted by $\mathcal{H}(G)$. From Principle A, it is also called the surface embedding graph of $G$.

Theorem 2 In $\mathcal{H}(G)$, there is a Hamilton path. Further, for any two vertices, $\mathcal{H}(G)$ has a Hamilton path with the two vertices as ends.

Proof Since a rotation at each vertex is a cyclic permutation(or in short a cycle) on the set of semi-edges with the vertex, an exchanger of layer segments is corresponding to a transposition on the set at a vertex.

Since any two cycles at a vertex $v$ can be transformed from one into another by $\rho(v)$ transpositions where $\rho(v)$ is the valency of $v$, i.e., the order of cycle(rotation), This enables us to do exchangers from the 1st layer on according to the order from left to right at one vertex to the other. Because of the finiteness, an associate polyhegon can always transformed into another by $|\mathcal{A}(G)|$ exchangers. From Theorem 1 with Principles $1-2$, the conclusion is done.

First, starting from a surface in $\mathcal{A}(G)$, by doing exchangers at each principle segments in one layer to another, a Hamilton path can always be found in considering Theorem 2 and Theorem 1. Then, a Hamilton path can be found on $\mathcal{H}(G)$.

Further, for chosen $S_{1}, S_{2} \in \mathcal{A}(G)=V(\mathcal{H}(G))$ adjective, starting from $S_{1}$, by doing exchangers avoid $S_{2}$ except the final step, on the basis of the strongly finite recursion principle, a Hamilton path between $S_{1}$ and $S_{2}$ can be obtained. In consequence, a Hamilton circuit can be found on $\mathcal{H}(G)$.

Corollary 1 In $\mathcal{H}(G)$, there exists a Hamilton circuit.

Theorem 2 tells us that the problem of determining the minimum, or maximum genus of graph $G$ has an algorithm in time linear on $\mathcal{H}(G)$.

## §3. Main Theorems

For a graph $G$, let $\mathcal{S}(G)$ be the the associate polehegons (or surfaces) of $G$, and $\mathbf{S}_{p}$ and $\mathbf{S}_{\tilde{q}}$, the subsets of, respectively, orientable and nonorientable polyhegons of genus $p \geqslant 0$ and $q \geqslant 1$.

Then, we have

$$
\mathcal{S}(G)=\sum_{p \geqslant 0} \mathbf{S}_{p}+\sum_{q \geqslant 1} \mathbf{S}_{\tilde{q}} .
$$

Theorem 3 A graph $G$ can be embedded on an orientable surface of genus $p$ if, and only if, $\mathcal{S}(G)$ has a polyhegon in $\mathbf{S}_{p}, p \geqslant 0$. Moreover, for an embedding of $G$, there exist a sequence of exchangers by which the corresponding polyhegon of the embedding can be transformed into one in $\mathbf{S}_{p}$.

Proof For an embedding of $G$ on an orientable surface of genus $p$, from Theorem 1 there is an associate polyhegon in $\mathbf{S}_{p}, p \geqslant 0$. This is the necessity of the first statement.

Conversely, given an associate polyhegen in $\mathbf{S}_{p}, p \geqslant 0$, from Theorems 1-2 with Principles A and B , an embedding of $G$ on an orientable surface of genus $p$ can be done. This is the sufficiency of the first statement.

The last statement of the theorem is directly seen from the proof of Theorem 2.
For an orientable embedding $\mu(G)$ of $G$, denote by $\widetilde{\mathbf{S}}_{\mu}$ the set of all nonorientable associate polyhegons induced from $\mu(G)$.

Theorem 4 A graph $G$ can be embedded on a nonorientable surface of genus $q(\geqslant 1)$ if, and only if, $\mathcal{S}(G)$ has a polyhegon in $\widetilde{\mathbf{S}}_{q}, q \geqslant 1$. Moreover, if $G$ has an embedding $\widetilde{\mu}$ on a nonorientable surface of genus $q$, then it can always be done from an orientable embedding $\mu$ arbitrarily given to another orientable embedding $\mu^{\prime}$ by a sequence of exchangers such that the associate polyhegon of $\widetilde{\mu}$ is in $\widetilde{\mathbf{S}}_{\mu^{\prime}}$.

Proof For an embedding of $G$ on a nonorientable surface of genus $q$, Theorem 1 and Principle B lead to that its associate polyhegon is in $\mathbf{S}_{q}, q \geqslant 1$. This is the necessity of the first statement.

Conversely, let $S_{\widetilde{q}}$ be an associate polyhegon of $G$ in $\widetilde{\mathbf{S}}_{q}, q \geqslant 1$. From Principles A and B, an embedding of $G$ on a nonorietable surface of genus $q$ can be found from $S_{\tilde{q}}$. This is the sufficiency of the first statement.

Since a nonorientable embedding of $G$ has exactly one under orientable embedding of $G$ by Principle A, Theorem 2 directly leads to the second statement.

## §4. Research Notes

A. Theorems 1 and 2 enable us to establish a procedure for finding all embeddings of a graph $G$ in linear space of the size of $G$ and in linear time of size of $\mathcal{H}(G)$. The implementation of this procedure on computers can be seen in [6].
B. In Theorems 3 and 4 , it is necessary to investigate a procedure to extract a sequence of transpositions considered for the corresponding purpose efficiently.
C. On the basis of the associate polyhegons, the recognition of operations from a polyhegon of genus $p$ to that of genus $p+k$ for given $k \geqslant 0$ have not yet be investigated. However, for the case $k=0$ the operations are just Operetions $0-2$ all topological that are shown in [1-3].
D. It looks worthful to investigate the associate polyhegon graph of a graph further for accessing the determination of the maximum(orientable) and minimum(orientable or nonorientable) genus of a graph.

## References

[1] Y.P. Liu, Advances in Combinatorial Maps(in Chinese), Northern Jiaotong University Press, Beijing, 2003.
[2] Y.P. Liu, Algebraic Principles of Maps(in Chinese), Higher Education Press, Beijing, 2006.
[3] Y.P. Liu, Theory of Polyhedra, Science Press, Beijing, 2008.
[4] Y.P. Liu, The nonorientable maximum genus of a graph(in Chinese with English abstract), Scintia Sinica, Special Issue on Math. I(1979), 191-201.
[5] Y.P. Liu, The maximum orientable genus of a graph, Scientia Sinica, Special Issue on Math. II(1979), 41-55.
[6] Wang Tao and Y.P. Liu, Implements of some new algorithms for combinatorial maps, $O R$ Trans., 12(2008), 2: 58-66.

# Plick Graphs with Crossing Number 1 

B.Basavanagoud<br>(Department of Mathematics, Karnatak University, Dharwad-580 003, India)<br>V.R.Kulli<br>(Department of Mathematics, Gulbarga University, Gulbarga-585 106, India)<br>E-mail: b.basavanagoud@gmail.com


#### Abstract

In this paper, we deduce a necessary and sufficient condition for graphs whose plick graphs have crossing number 1 . We also obtain a necessary and sufficient condition for plick graphs to have crossing number 1 in terms of forbidden subgraphs.


Key Words: Smarandache $\mathscr{P}$-drawing, drawing, line graph, plick graph, crossing number.
AMS(2010): 05C10

## §1. Introduction

All graphs considered here are finite, undirected and without loops or multiple edges. We refer the terminology of [2]. For any graph $G, L(G)$ denote the line graph of $G$.

A Smarandache $\mathscr{P}$-drawing of a graph $G$ for a graphical property $\mathscr{P}$ is such a good drawing of $G$ on the plane with minimal intersections for its each subgraph $H \in \mathscr{P}$. A Smarandache $\mathscr{P}$-drawing is said to be optimal if $\mathscr{P}=G$ and it minimizes the number of crossings. A graph is planar if it can be drawn in the plane or on the sphere in such a way that no two of its edges intersect. The crossing number $\operatorname{cr}(G)$ of a graph $G$ is the least number of intersections of pairs of edges in any embedding of $G$ in the plane. Obviously, $G$ is planar if and only if $\operatorname{cr}(G)=0$. It is implicit that the edges in a drawing are Jordan $\operatorname{arcs}$ (hence, non-selfintersecting), and it is easy to see that a drawing with the minimum number of crossings(an optimal drawing) must be good drawing, that is, each two edges have at most one vertex in common, which is either a common end-vertex or a crossing. Theta is the result of adding a new edge to a cycle and it is denoted by $\theta$. The corona $G^{+}$of a graph $G$ is obtained from $G$ by attaching a path of length 1 to every vertex of $G$.

The plick graph $P(G)$ of a graph $G$ is obtained from the line graph by adding a new vertex corresponding to each block of the original graph and joining this vertex to the vertices of the line graph which correspond to the edges of the block of the original graph(see[4]).

The following will be useful in the proof of our results.

[^2]Theorem $\mathbf{A}([5])$ The line graph of a planar graph $G$ is planar if and only if $\Delta(G) \leq 4$ and every vertex of degree 4 is a cut-vertex.

Theorem $\mathbf{B}([3])$ Let $G$ be a nonplanar graph. Then $\operatorname{cr}(L(G))=1$ if and only if the following conditions hold:
(1) $\operatorname{cr}(G)=1$;
(2) $\Delta(G) \leq 4$, and every vertex of degree 4 is a cut-vertex of $G$;
(3) There exists a drawing of $G$ in the plane with exactly one crossing in which each crossed edge is incident with a vertex of degree 2.

Theorem $\mathbf{C}([3])$ The line graph of a planar graph $G$ has crossing number one if and only if (1) or (2) holds:
(1) $\Delta(G)=4$ and there is a unique non-cut-vertex of degree 4;
(2) $\Delta(G)=5$, every vertex of degree 4 is a cut-vertex, there is a unique vertex of degree 5 and it has at most 3 incident edges in any block.

Theorem $\mathbf{D}([4])$ The plick graph $P(G)$ of a graph $G$ is planar if and only if $G$ satisfies the following conditions:
(1) $\Delta(G) \leq 4$, and
(2) every block of $G$ is either a cycle or a $K_{2}$.

Theorem $\mathbf{E}([1])$ A graph has a planar ilne graph if and only if it has no subgraph homeomorphic to $K_{3,3}, K_{1,5}, P_{4}+K_{1}$ or $K_{2}+\bar{K}_{3}$.

Remark 1([4]) For any graph, $L(G)$ is a subgraph of $P(G)$.

## §2. Results

The following theorem supports the main theorem.

Theorem 1 Let $x$ be any edge of $K_{4}$. If $G$ is homeomorphic to $K_{4}-x$, then $\operatorname{cr}(P(G))=1$.

Proof We prove the theorem first for $G=\left(K_{4}-x\right)$. One can see that the graph $P\left(K_{4}-x\right)$ has 6 vertices and 13 edges. But a planar graph with 6 vertices has at most 12 edges. This shows that $P\left(K_{4}-x\right)$ has crossing number at least 1. Figure 1, being drawing of $P\left(K_{4}-x\right)$ concludes that $\operatorname{cr}\left(P\left(K_{4}-x\right)\right)=1$. Suppose now $G$ is the graph as in the statement. Referring to Figure 1, it is immediate to see that $\operatorname{cr}\left(P\left(K_{4}-x\right)\right)=1$.


Figure 1
The following theorem gives a necessary and sufficient condition for graphs whose plick graphs have crossing number 1 .

Theorem 2 A graph $G$ has a plick graph with crossing number 1 if and only if $G$ is planar and one of the following holds:
(1) $\Delta(G)=3, G$ has exactly two non-cut-vertices of degree 3 and they are adjacent.
(2) $\Delta(G)=4$, every vertex of degree 4 is a cut-vertex of $G$, there exists exactly one theta as a block in $G$ such that at least one vertex of theta is a non-cut- vertex of degree 2 or 3 and every other block of $G$ is either a cycle or a $K_{2}$.
(3) $\Delta(G)=5, G$ has a unique cut-vertex of degree 5 and every block of $G$ is either a cycle or a $K_{2}$.

Proof Suppose $P(G)$ has crossing number one. Then by Remark 1, and Theorem B, $G$ is planar. By Theorem $\mathrm{D}, \Delta(G) \leq 4$, then at least one block of $G$ is neither a cycle nor a $K_{2}$.

Suppose $\Delta(G) \leq 6$. Then $K_{1,6}$ is a subgraph of $G$. Clearly $L\left(K_{1,6}\right)=K_{6}$. It is known that $\operatorname{cr}\left(K_{6}\right)=3$. By Remark 1, $K_{6}$ is a subgraph of $P(G)$ and hence $\operatorname{cr}(P(G))>1$, a contradiction. This implies that $\Delta(G) \leq 5$. If $\Delta(G) \leq 2$, then $P(G)$ is planar, again a contradiction. Thus $\Delta(G)=3$ or 4 or 5 .

We now consider the following cases:
Case 1. Suppose $\Delta(G)=3$. Then by Theorem D and since $\operatorname{cr}(P(G))=1, G$ has a non-cut-vertex of degree 3 . Clearly $G$ contains a subgraph homeomorphic to $K_{4}-x$, so that there exist at least two non-cut-vertex of degree3. More precisely, there is an even number, say $2 n$, of non-cut-vertex of degree 3 . Now suppose $G$ has at least two diagonal edges. Then there are two subcases to consider depending on whether 2 diagonal edges exist in one cycle or in two different edge disjoint cycles.

Subcase 1.1 If two diagonal edges exist in one cycle of $G$. Then $G$ has a subgraph homeomorphic from $K_{4}$. The graph $P\left(K_{4}\right)$ has 7 vertices and 18 edges. It is known that a planar graph with 7 vertices has at most 15 edges. This shows that $P\left(K_{4}\right)$ must have crossing number exceeding 1 and hence $P(G)$ has crossing number greater than 1 , a contradiction.

Subcase 1.2 If two diagonal edges exist in two different edge-disjoint cycles of $G$. Then by

Theorem 1, we see that for every subgraph of $G$ homeomorphic to $K_{4}-x$, there corresponds at least one crossing of $G$. Hence $P(G)$ has at least 2 crossings, a contradiction.

Hence $G$ has exactly two non-cut-vertices of degree 3 and every other vertex of degree 3 is a cut-vertex.

Suppose a graph $G$ has two non-cut-vertices of degree 3 and they are not adjacent. Then $G$ contains a subgraph homeomorphic to $K_{2,3}$. On drawing $P\left(K_{2,3}\right)$ in a plane one can see that $\operatorname{cr}\left(P\left(K_{2,3}\right)\right)=2$. Since $P\left(K_{2,3}\right)$ is a subgraph of $P(G), P(G)$ has crossing number exceeding 1, a contradiction(see Figure 2).

Therefore, we conclude that $G$ contains exactly two non-cut-vertices of degree 3 and these are adjacent. This proves (1).


Figure 2
Case 2. Assume $\Delta(G)=4$. We show first that every vertex of degree 4 is a cut-vertex. On the contrary suppose that $G$ has non-cut-vertex $v$ of degree 4 . Then by Theorem C, $\operatorname{cr}(L(G)) \geq 1$. The vertex $u_{1}$ in $P(G)$ corresponding to the block which contains a non-cut-vertex of degree 4 is adjacent to every vertex of $L(G)$. We obtain the drawing of $P(G)$ with 3 crossings.

Assume now $G$ has at least two blocks each of which is a theta. By Theorem 1 and case 1 of this theorem, we see that for every subgraph of $G$ homeomorphic to $K_{4}-x$, there correspond to at least 2 crossings of $G$, a contradiction.

Suppose there exists exactly one theta $S$ as a block in $G$ such that none of its vertices is a non-cut-vertex of degree 2 or 3 . Assume all vertices of theta $S$ have degree 4 in $G$. Then by Theorem A, $L(S)$ is planar. Let $v_{1}$ be the vertex of $L(G)$ corresponding to the chord of a cycle $C$ of theta. The vertex $w_{1}$ in $P(G)$ corresponding to the block theta $S$ is adjacent to every vertex of $L(C)$ without crossings. In $P(G)-v_{1} w_{1}$, the vertex $w_{1}$ is adjacent to every vertex of $L(S)-v_{1}$ without crossings. By the definition of $P(G)$, the vertices $v_{1}$ and $w_{1}$ are adjacent in $P(G)$. The edge $v_{1} w_{1}$ crosses at least two edges of $L(G)$. On drawing of $P(G)$ in the plane, it has at least two crossings, a contradiction. This proves that $\Delta(G)=4$, there exist exactly one theta as a block in $G$ such that at least one vertex of theta is either a non-cut-vertex of degree 2 or 3 .

Suppose every block of $G$ different from theta block is neither a cycle nor a $K_{2}$. It implies that $G$ has a block which is a subgraph homeomorphic to $K_{4}-x$. By Cases 1 and 2 of this theorem, we see that for every subgraph of $G$ homeomorphic to $K_{4}-x$, there corresponds at
least one crossing of $G$. Hence $P(G)$ has at least 2 crossings, a contradiction.


Figure 3


## Figure 4

Case 3. Assume $\Delta(G)=5$. Suppose $G$ has at least two vertices of degree 5 . Then by Theorem C, $L(G)$ has crossing number at least 2. By Remark $1, \operatorname{cr}(P(G)) \geq 2$, which is a contradiction. Thus $G$ has a unique vertex of degree 5 .

Suppose $G$ has a vertex $v$ of degree 5 and at least one block of $G$ is neither a cycle nor a $K_{2}$. Then some block of $G$ has a subgraph homeomorphic to $K_{4}-x$. By Case 1 of this theorem $\operatorname{cr}\left(P\left(K_{4}-x\right)\right) \geq 1$ and the 5 edges incident to $v$ form $K_{5}$ as a subgraph in $P(G)$. Hence $\operatorname{cr}\left(P\left(K_{4}-x\right)\right) \geq 2$, a contradiction.

Conversely, suppose $G$ is a planar graph satisfying (1) or (2) or (3). Then by Theorem $\mathrm{D}, P(G)$ has crossing number at least 1 . We now show that its crossing number is at most 1 . First suppose (1) holds. Then $G$ has exactly one block, say $H$, homeomorphic to $K_{4}-x$ which contains 2 adjacent non-cut-vertices of degree 3. By Theorem $1, \operatorname{cr}(P(H))=1$. By Theorem D, all other remaining blocks of $G$ have a planar plick graph. Hence $P(G)$ has crossing number 1.

Assume (2) holds. Let $u$ be a cut-vertex of degree 4. The vertex $u$ has a non-cut-vertex of degree 3 in a block for otherwise, $G$ does contain a subgraph homeomorphic to $K_{4}-x$ which is impossible. By virtue of Theorem 1, for a non-cut-vertex of $G$ of degree 3, there corresponds one crossing in $P(G)$. However $P(G)$ can not have more than one crossing since the removal of any edge in a block containing $u$, yields a graph $H$ such that $P(H)$ is planar by Theorem D . It follows easily that $P(G)$ has crossing number 1 .

Suppose (3) holds. The edges at the vertex $v$ of the degree 5 can be split into sets of sizes

2 and 3 so that no edges in different sets are in the same block. Transform $G$ to $G^{\prime}$ as in Figure 3. Then $P\left(G^{\prime}\right)$ is again planar. Thus $P(G)$ can be drawn with only one crossing as shown in Figure 4.

## §3. Forbidden Subgraphs

By using Theorem 2, we now characterize graphs whose plick graphs have crossing number 1 in terms of forbidden subgraphs.

Theorem 3 The plick graph of a connected graph $G$ has crossing number 1 if and only if $G$ has no subgraphs homeomorphic from any one of the graphs of Figure 5 or $G$ has subgraph $\theta^{+}$ such that none of the vertices of theta have non-cut-vertices of degree 2 or 3.

$G_{3}:$


Figure 5

Proof Suppose $G$ has a plick graph with crossing number one. We now show that all graphs homeomorphic from any one of the graphs of Figure 5 or a subgraph $\theta^{+}$such that none of the vertices of theta have non-cut-vertices of degree 2 or 3 , have no plick graph with crossing number one. This result follows from Theorem 2, since graphs homeomorphic from $G_{1}, G_{2}$ or $G_{3}$ have more than two non-cut-vertices of degree three. Graphs homeomorphic from $G_{4}$ have two non-cut-vertices of degree 3 which are not adjacent. Graphs homeomorphic from
$G_{5}$ have a vertex of degree 4 which is a non-cut-vertex. Graphs homeomorphic from $G_{6}$ have more than one theta. $\theta^{+}$has exactly one block which is a theta and none of its vertices have non-cut-vertices of degree 2 or 3 . Graphs homeomorphic from $G_{7}$ have $\Delta\left(G_{7}\right)>5$. Graphs homeomorphic from $G_{8}$ or $G_{9}$ have two or more vertices of degree 5. Graphs homeomorphic to $G_{10}$ or $G_{11}$ have a block which is neither a cycle nor a $K_{2}$.

Conversely, suppose $G$ is a graph which does not contain a subgraph homeomorphic from any one of the graphs of Figure 5 or $G$ has exactly one subgraph theta as a block such that none of the vertices of theta have non-cut-vertices of degree 2 or 3 . First we prove condition (1) of Theorem 2. Suppose $G$ contains more than two non-cut-vertices of degree 3. Then it is easy to see that $G$ is a planar graph with at least 2 diagonal edges. Now consider 2 cases depending on whether the 2 diagonal edges exist in one block or in two different blocks.

Case 1. Suppose two diagonal edges exist in one block of $G$, then $G$ has a subgraph homeomorphic from $G_{1}$ or $G_{2}$.

Case 2. Suppose two diagonal edges exist in two different blocks of $G$, then $G$ has a subgraph homeomorphic from $G_{3}$.

In each case we have a contradiction. Hence $G$ has at most two non-cut-vertices of degree 3. Suppose $G$ has exactly two nonadjacent non-cut-vertices of degree 3. Then there exist 3 disjoint paths between these two non-cut-vertices of degree 3. Clearly $G$ contains a subgraph homeomorphic from $G_{4}$, a contradiction. Thus $G$ has exactly two adjacent non-cut-vertices of degree 3.

Since $G$ does not contain a subgraph homeomorphic from $G_{7}$ i.e, $K_{1,6}, \Delta(G) \leq 5$. Also since $\Delta(G) \geq 4$, if it follows that $\Delta(G)=4$ or 5 .

Suppose $G$ has a vertex $v$ of degree 4. We prove that $v$ is a cut-vertex. If not, let $a, b, c$ and $d$ be the vertices of $G$ adjacent to $v$. Then there exist paths between every pair of vertices of $a, b, c$ and $d$ not containing $v$. Then it is proved in Theorem $\mathrm{E}, G$ has a subgraph homeomorphic from $G_{5}$, this is a contradiction. Thus $v$ is a cut-vertex and every vertex of degree 4 is a cut-vertex.

Suppose that a cut-vertex of degree 4 lies on two blocks, each of which is a theta. Then $G$ has a subgraph homeomorphic from $G_{6}$. This is a contradiction. $G$ has exactly one block which is a theta such that at least one vertex of theta is either a non-cut-vertex of degree 2 or 3 , for otherwise a forbidden subgraph has exactly one theta as a block such that none of the vertices of theta have non-cut-vertices of degree 2 or 3 would appear in $G$.

Suppose $G$ has two vertices $v_{1}$ and $v_{2}$ of degree 5 . Since $G$ is a connected, $v_{1}$ and $v_{2}$ are connected by a path $P$ and let $\left(v_{1}, a_{i}\right)$ and $\left(v_{2}, b_{j}\right), i, j=1,2,3,4$, be edges of $G$. We consider the following possibilities.

If $a_{i} \neq b_{j}$ for $i, j=1,2,3,4$, then $G$ contains a subgraph homeomorphic from $G_{8}$, a contradiction.

If there exists a path between a vertex of $a_{i}$ and a vertex of $b_{j}$, then $G$ has a subgraph homeomorphic from $G_{9}$, a contradiction.

If $a_{i}=b_{j}$, for $i, j=1,2$, then clearly $G$ contains a subgraph homeomorphic from $G_{10}$, a contradiction.

This proves that $G$ has exactly one vertex $v$ of degree 5 .

Suppose $G$ has a vertex $v$ of degree 5 . We show that $v$ is a cut-vertex. If possible let us assume that $G$ has a non-cut-vertex of degree 5. In this case Greenwell and Hemminger showed in [1] that $G$ must contain a subgraph homeomorphic from $G_{5}$, a contradiction.

Suppose $G$ has a unique cut-vertex $v$ of degree 5 and it lies on blocks, one block which is neither a cycle nor a $K_{2}$. Then $G$ contains a subgraph homeomorphic from $G_{10}$ or $G_{11}$.

Thus Theorem 2 implies that $G$ has a plick graph with crossing number one.

## References

[1] D.L.Greenwell and R.L. Hemminger, Forbidden subgraphs for graphs with planar line graphs, Discrete Math., 2(1972), 31-34.
[2] F.Harary, Graph Theory, Addison-Wesley, Reading, Mass, (1969).
[3] S.Jendroĺ, M.Klešč, On graphs whose line graphs have crossing number one, J.Graph Theory, 37(2001), 181-188.
[4] V.R.Kulli and B.Basavanagoud, Characterizations of planar plick graphs, Discussiones Mathematicae, Graph Theory, 24(2004), 41-45.
[5] J.Sėdláček, Some properties of interchange graphs, in: Theory of graphs and its applications, M.Fiedler, ed.(Academic Press, New York, 1962), 145-150.

# Effects of Foldings on Free Product of Fundamental Groups 

M.El-Ghoul, A. E.El-Ahmady, H.Rafat and M.Abu-Saleem

(Mathematics Department, Faculty of Science, Tanta University, Tanta, Egypt)

E-mail: hishamrafat2005@yahoo.com


#### Abstract

In this paper, we will introduce free fundamental groups of some types of manifolds. Some types of conditional foldings restricted on the elements on free group and their fundamental groups are deduced. Also, the fundamental group of the limit of foldings on a wedge sum of two manifolds is obtained. Theorems governing these relations will be achieved.


Key Words: Manifolds, Folding, fundamental group, Free group
AMS(2010): 51H20, 57N10, 57M05,14F35,20F34

## §1. Introduction

In this article the concept of foldings will be discussed from viewpoint of algebra. The effect of foldings on the manifold $M$ or on a finite number of product manifolds $M_{1} x M_{2} x \ldots x M_{n}$ on the fundamental group $\pi_{1}(M)$ and $\pi_{1}\left(M_{1} x M_{2} x \cdots x M_{n}\right)$ will be presented. The folding of a manifold was, firstly introduced by Robertson 1977 [14]. More studies on the folding of many types of manifolds were studied in [2-4 and 6-9]. The unfolding of a manifold introduced in [5]. Some application of the folding of a manifold discussed in [1]. The fundamental groups of some types of a manifold are discussed in [10-13].

## §2. Definitions

1. The set of homotopy classes of loops based at the point $x_{\circ}$ with the product operation $[f][g]=[f \cdot g]$ is called the fundamental group and denoted by $\pi_{1}\left(X, x_{\circ}\right)[11]$.
2. Let $M$ and $N$ be two smooth manifolds of dimension $m$ and $n$ respectively. A map $f: M \rightarrow N$ is said to be an isometric folding of $M$ into $N$ if for every piecewise geodesic path $\gamma: I \rightarrow M$ the induced path $f \circ \gamma: I \rightarrow N$ is piecewise geodesic and of the same length as $\gamma$ [14]. If $f$ does not preserve length it is called topological folding [9].
3. Let $M$ and $N$ be two smooth manifolds of the same dimension. A map $g: M \rightarrow N$ is said to be unfolding of $M$ into $N$ if every piecewise geodesic path $\gamma: I \rightarrow M$, the induced path $g \circ \gamma: I \rightarrow N$ is piecewise geodesic with length greater than $\gamma[5]$.

[^3]4. Given spaces $X$ and $Y$ with chosen points $x_{o} \in X$ and $y_{o} \in Y$, then the wedge sum $X \vee Y$ is the quotient of the disjoint union $X \cup Y$ obtained identifying $x_{o}$ and $y_{o}$ to a single point [11].
5. Let S be an arbitrary set. A free group on the set S is a group $F$ together with a function $\phi: S \rightarrow F$ such that the following condition holds: For any function $\psi: S \rightarrow H$, there exist a unique homomorphism $f: F \rightarrow H$ such that $f \circ \phi=\psi[12]$.

## §3. Main Results

Paving the stage to this paper, we then introduce the following
(1) $\pi_{1}(T)=\left\{\left(\left[\alpha_{1}\right]^{k},\left[\beta_{1}\right]^{m}\right),\left(\left[\alpha_{2}\right]^{k},\left[\beta_{2}\right]^{m}\right), \ldots .,\left(\left[\alpha_{n}\right]^{k},\left[\beta_{n}\right]^{m}\right) ;\left[\alpha_{i}\right],\left[\beta_{i}\right] \in \pi_{1}\left(S^{1}\right), k, m \in Z, k \neq\right.$ $0, m \neq 0, i=1,2, \ldots, n\}$
(2) $\pi_{1}(T) \bmod (k, m)=\left\{\left(\left[\alpha_{1}\right],\left[\beta_{1}\right]\right),\left(\left[\alpha_{2}\right],\left[\beta_{2}\right]\right), \ldots,,\left(\left[\alpha_{n}\right],\left[\beta_{n}\right]\right):\left[\alpha_{i}\right]^{k}=1,\left[\beta_{i}\right]^{m}=1\left[\alpha_{i}\right],\left[\beta_{i}\right] \in\right.$ $\left.\pi_{1}\left(S^{1}\right), k, m \in Z^{+}, k \neq 0, m \neq 0, i=1,2, \ldots, n\right\}$.

Where, $\pi_{1}\left(S^{1}\right)$ is a fundamental group of the circle,$T$ is the torus $[\alpha]^{n}=\underbrace{[\alpha] \times[\alpha] \times \ldots \times[\alpha]}_{n-\text { terms }}$, and $T^{n}=\underbrace{T \times T \times \ldots \times T}_{n-\text { terms }}$.

Let $\pi_{1}\left(S_{1}^{1}\right), \pi_{1}\left(S_{2}^{1}\right)$ be two fundamental groups. Then the free product of $\pi_{1}\left(S_{1}^{1}\right), \pi_{1}\left(S_{2}^{1}\right)$ is the group $\pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ consisting of all reduced words $a_{1} a_{2} a_{3} \ldots . a_{m}$ of an arbitrary finite length $m \geq 0$ such that $a_{i} \in \pi_{1}\left(S_{1}^{1}\right)$ or $a_{i} \in \pi_{1}\left(S_{2}^{1}\right), i=1,2, \ldots, m$, then we can represent the elements $a_{i}$ as of the forms $a_{i}=[\alpha]^{n_{i}}$ or $a_{i}=[\beta]^{n_{i}}$ where $n_{i} \in Z, n_{i} \neq 0$ and $\alpha, \beta$ are two loops that goes once a round $S_{1}^{1}$, $S_{2}^{1}$ respectively. Also, if $F: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}$ is a folding, then the induced folding $\overline{F:} \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ has the following forms:

$$
\begin{aligned}
& \left.\left.\overline{F( } \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)\right)=\overline{F( } \pi_{1}\left(S_{1}^{1}\right)\right) * \pi_{1}\left(S_{2}^{1}\right) \\
& \left.\left.\overline{F( } \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)\right)=\pi_{1}\left(S_{1}^{1}\right) * \overline{F( } \pi_{1}\left(S_{2}^{1}\right)\right) \\
& \left.\left.\left.\overline{F( } \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)\right)=\overline{F( } \pi_{1}\left(S_{1}^{1}\right)\right) * \overline{F( } \pi_{1}\left(S_{2}^{1}\right)\right)
\end{aligned}
$$

Theorem 3.1 If $F_{i}: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1} \quad, i=1,2$ are two types of of foldings, where $F_{I}\left(S_{j}^{1}\right)=. S_{j}^{1}, j=1,2$, then there are induced foldings $\overline{F_{i}}: \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ such that $\overline{F_{i}}\left(\pi_{1}\left(S_{1}^{1}\right)\right) * \pi_{1}\left(S_{2}^{1}\right) \approx Z$.

Proof First, let $F_{1}: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}$ is folding such that $F_{1}\left(S_{1}^{1}\right)=S_{1}^{1}, F_{1}\left(S_{2}^{1}\right)=S_{1}^{1}$ as in Fig.1. Then we can express each element $g=a_{1} a_{2} a_{3} \ldots . a_{m}, m \geq 1$ of $\pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ in the following forms

$$
\begin{aligned}
& {[\alpha]^{n_{1}}[\beta]^{n_{2}}[\alpha]^{n_{3}} \ldots[\alpha]^{n_{m-1}}[\beta]^{n_{m}},[\alpha]^{n_{1}}[\beta]^{n_{2}}[\alpha]^{n_{3}} \cdots[\beta]^{n_{m-1}}[\alpha]^{n_{m}},} \\
& {[\beta]^{n_{1}}[\alpha]^{n_{2}}[\beta]^{n_{3}} \cdots[\beta]^{n_{m-1}}[\alpha]^{n_{m}}, \text { or }[\beta]^{n_{1}}[\alpha]^{n_{2}}[\beta]^{n_{3}} \ldots[\alpha]^{n_{m-1}}[\beta]^{n_{m}},}
\end{aligned}
$$

where $n_{1}, n_{2}, \cdots, n_{m}$ are nonzero integers and $[\alpha]^{n_{k}} \in \pi_{1}\left(S_{1}^{1}\right),[\beta]^{n_{k}} \in \pi_{1}\left(S_{2}^{1}\right), k=1,2, . . m$.

Then, the induced folding of the element $g$ is

$$
\begin{aligned}
& \overline{F_{1}}(g)=\overline{F_{1}}\left([\alpha]^{n_{1}}\right) \overline{F_{1}}\left([\beta]^{n_{2}}\right) \overline{F_{1}}\left([\alpha]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\alpha]^{n_{m-1}}\right) \overline{F_{1}}\left([\beta]^{n_{m}}\right), \\
& \overline{F_{1}}\left([\alpha]^{n_{1}}\right) \overline{F_{1}}\left([\beta]^{n_{2}}\right) \overline{F_{1}}\left([\alpha]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\beta]^{n_{m-1}}\right) \overline{F_{1}}\left([\alpha]^{n_{m}}\right), \\
& \overline{F_{1}}\left([\beta]^{n_{1}}\right) \overline{F_{1}}\left([\alpha]^{n_{2}}\right) \overline{F_{1}}\left([\beta]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\beta]^{n_{m-1}}\right) \overline{F_{1}}\left([\alpha]^{n_{m}}\right), \\
& \overline{F_{1}}\left([\beta]^{n_{1}}\right) \overline{F_{1}}\left([\alpha]^{n_{2}}\right) \overline{F_{1}}\left([\beta]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\alpha]^{n_{m-1}}\right) \overline{F_{1}}\left([\beta]^{n_{m}}\right) .
\end{aligned}
$$

Since $\overline{F_{1}}\left([\alpha]^{n_{k}}\right)=[\alpha]^{n_{k}}, \overline{F_{1}}\left([\beta]^{n_{k}}\right)=[\alpha]^{n_{k}}$ it follows that $\overline{F_{1}}\left(a_{1} a_{2} a_{3} \ldots a\right)=[\alpha]^{\left(n_{1}+n_{2}+\cdots+n_{m}\right)}$. Hence, there is an induced folding $\overline{F_{i}}: \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ such that $\overline{F_{i}}\left(\pi_{1}\left(S_{1}^{1}\right) *\right.$ $\left.\pi_{1}\left(S_{2}^{1}\right)\right)=\pi_{1}\left(S_{1}^{1}\right)$, and so $\overline{F_{i}}\left(\pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)\right) \approx Z$. Similarly, if $F_{2}: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}$ is folding, such that $F_{2}\left(S_{1}^{1}\right)=S_{2}^{1}, F_{2}\left(S_{2}^{1}\right)=S_{2}^{1}$, then there is an induced folding $\overline{F_{2}}: \pi_{1}\left(S_{1}^{1}\right) *$ $\pi_{1}\left(S_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ such that $\overline{F_{2}}\left(\pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)\right) \approx Z$.


Fig. 1

Theorem 3.2 If $F_{i}: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}, i=1,2$ are two types of foldings such that $F_{i}\left(S_{j}^{1}\right)=$ $S_{i}^{1}, j=1,2$. Then, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)$ is isomorphic to $Z$.

Proof Let $F_{i}\left(S_{j}^{1}\right)=S_{i}^{1}$ then $\lim _{n \rightarrow \infty} F_{i_{n}}\left(S_{1}^{1} \vee S_{2}^{1}\right)=S_{i}^{1}$ as in Fig.2. Thus, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)$ $=S_{i}^{1}$, Therefore $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)$ is isomorphic to $Z$.


Fig. 2
Theorem 3.3 Let $F: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1} \quad$ be a folding, where $F\left(S_{i}^{1}\right) \neq S_{i}^{1}, i=1$, 2. Then there is an induced folding $\bar{F}: \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ such that $\left.\bar{F} \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)\right)=0$.

Proof Let $F: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}$ be a folding such that $F\left(S_{1}^{1}\right) \neq S_{1}^{1}, F\left(S_{i}^{1}\right) \neq S_{i}^{1}$ as in Fig. (3).Then, we can express each element $g=a_{1} a_{2} a_{3} \ldots a_{m}, m \geq 1$ of $\pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ in the following forms:

$$
\begin{array}{ll}
{[\alpha]^{n_{1}}[\beta]^{n_{2}}[\alpha]^{n_{3}} \cdots[\alpha]^{n_{m-1}}[\beta]^{n_{m}},} & {[\alpha]^{n_{1}}[\beta]^{n_{2}}[\alpha]^{n_{3}} \cdots[\beta]^{n_{m-1}}[\alpha]^{n_{m}}} \\
{[\beta]^{n_{1}}[\alpha]^{n_{2}}[\beta]^{n_{3}} \cdots[\beta]^{n_{m-1}}[\alpha]^{n_{m}},} & {[\beta]^{n_{1}}[\alpha]^{n_{2}}[\beta]^{n_{3}} \cdots[\alpha]^{n_{m-1}}[\beta]_{m}^{n_{m}}}
\end{array}
$$

where $n_{1}, n_{2}, \cdots, n_{m}$ are nonzero integers and $[\alpha]^{n_{k}} \in \pi_{1}\left(S_{1}^{1}\right),[\beta]^{n_{k}} \in \pi_{1}\left(S_{2}^{1}\right), k=1,2, \cdots, m$.

Then the induced folding of the element $g$ is

$$
\begin{aligned}
& \overline{F_{1}}(g)=\overline{F_{1}}\left([\alpha]^{n_{1}}\right) \overline{F_{1}}\left([\beta]^{n_{2}}\right) \overline{F_{1}}\left([\alpha]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\alpha]^{n_{m-1}}\right) \overline{F_{1}}\left([\beta]^{n_{m}}\right) \\
& =[\alpha]^{n_{1}}[\beta]^{n_{2}}[\alpha]^{n_{3}} \cdots[\alpha]^{n_{m-1}}[\beta]^{n_{m}}, \\
& \overline{F_{1}}\left([\alpha]^{n_{1}}\right) \overline{F_{1}}\left([\beta]^{n_{2}}\right) \overline{F_{1}}\left([\alpha]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\beta]^{n_{m-1}}\right) \overline{F_{1}}\left([\alpha]^{n_{m}}\right) \\
& =[\alpha]^{n_{1}}[\beta]^{n_{2}}[\alpha]^{n_{3}} \cdots[\beta]^{n_{m-1}}[\alpha]^{n_{m}}, \\
& \overline{F_{1}}\left([\beta]^{n_{1}}\right) \overline{F_{1}}\left([\alpha]^{n_{2}}\right) \overline{F_{1}}\left([\beta]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\beta]^{n_{m-1}}\right) \overline{F_{1}}\left([\alpha]^{n_{m}}\right) \\
& =[\beta]^{n_{1}}[\alpha]^{n_{2}}[\beta]^{n_{3}} \cdots[\beta]^{n_{m-1}}[\alpha]^{n_{m}}, \\
& \overline{F_{1}}\left([\beta]^{n_{1}}\right) \overline{F_{1}}\left([\alpha]^{n_{2}}\right) \overline{F_{1}}\left([\beta]^{n_{3}}\right) \cdots \overline{F_{1}}\left([\alpha]^{n_{m-1}}\right) \overline{F_{1}}\left([\beta]^{n_{m}}\right) \\
& =[\beta]^{n_{1}}[\alpha]^{n_{2}}[\beta]^{n_{3}} \cdots[\alpha]^{n_{m-1}}[\beta]^{n_{m}} .
\end{aligned}
$$

It follows from $\left[\begin{array}{l}\hat{\alpha}],[\hat{\beta}] \longrightarrow \text { identity element, that there is an induced folding } \overline{F:} \pi_{1}\left(S_{1}^{1}\right) *\end{array}\right.$ $\pi_{1}\left(S_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ such that $\left.\overline{F( } \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)\right)=0$.


Fig. 3
Corollary 1 If $F_{i}: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}, i=1,2$ are two types of foldings such that

$$
F_{i}\left(S_{i}^{1}\right)=S_{i}^{1}, F_{j}\left(S_{i}^{1}\right) \neq S_{i}^{1}, j=1,2, i \neq j
$$

Then there are induced foldings $\overline{F_{i}}: \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$ such that $\overline{F_{i}}\left(\pi_{1}\left(S_{1}^{1}\right) *\right.$ $\left.\pi_{1}\left(S_{2}^{1}\right)\right) \approx Z$.

Theorem 4 If $F: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}$ is a folding such that $F\left(S_{i}^{1}\right) \neq S_{i}^{1}, i=1,2$. Then,

$$
\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)
$$

is the identity group.
Proof If $F\left(S_{i}^{1}\right) \neq S_{i}^{1}, i=1,2$ then $\lim _{n \rightarrow \infty} F_{n}\left(S_{1}^{1} \vee S_{2}^{1}\right)$ is a point as in Fig.4, and so $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)$ is the fundamental group of a point. Therefore, we get that $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(S_{1}^{1} \vee\right.\right.$ $\left.\left.S_{2}^{1}\right)\right)=0$.


Fig. 4
Theorem 5 If $F_{i}: S_{1}^{1} \vee S_{2}^{1} \longrightarrow S_{1}^{1} \vee S_{2}^{1}, i=1,2$ are two types of foldings such that $F_{i}\left(S_{i}^{1}\right)=S_{i}^{1}, F_{j}\left(S_{i}^{1}\right) \neq S_{i}^{1}, j=1,2, i \neq j$. Then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)$ is isomorphic to $Z$.

Proof It follows from $F_{i}\left(S_{i}^{1}\right)=S_{i}^{1}, F_{j}\left(S_{i}^{1}\right) \neq S_{i}^{1}, j=1,2, i \neq j$.that the limit of one circle is a circle and the limit of the other circle is a point, so $\left.\lim _{n \rightarrow \infty} F_{n}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)=S_{i}^{1}$ as in Fig.5. Thus, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)=\pi_{1}\left(S_{i}^{1}\right)$. Therefore $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(S_{1}^{1} \vee S_{2}^{1}\right)\right)$ is isomorphic to $Z$. $\square$


Fig. 5
Now, we will generalize the above concepts for the tours Consider $\pi_{1}\left(T_{1}^{1}\right), \pi_{1}\left(T_{2}^{1}\right)$, are two fundamental groups. Then, the free product of $\pi_{1}\left(T_{1}^{1}\right), \pi_{1}\left(T_{2}^{1}\right)$, is the group $\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$ consisting of all reduced words of $a_{1} a_{2} a_{3} \ldots a_{m}$ of an arbitrary finite length $m \geq 0$ such that
$a_{i} \in \pi_{1}\left(T_{1}^{1}\right)$ or $a_{i} \in \pi_{1}\left(T_{2}^{1}\right)$ and so, we can represent the elements $a_{i}$ as of the forms $a_{i}=$ $\left(\left[\alpha_{1}\right]^{n_{i}},\left[\beta_{1}\right]^{k_{i}}\right)$ or $a_{i}=\left(\left[\alpha_{2}\right]^{n_{i}},\left[\beta_{2}\right]^{k_{i}}\right)$ where $n_{i}, k_{i} \in Z, n_{i} \neq 0, k_{i} \neq 0$ where $\left(\left[\alpha_{1}\right]^{n_{i}},\left[\beta_{1}\right]^{k_{i}}\right) \in$ $\pi_{1}\left(T_{1}^{1}\right),\left(\left[\alpha_{2}\right]^{n_{i}},\left[\beta_{2}\right]^{k_{i}}\right) \in \pi_{1}\left(T_{2}^{1}\right)$ and $\alpha_{j}, \beta_{j}$ are loops that goes once a round the generators of $T_{j}$ for $j=1,2$.Then if $F: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}$ is a folding, then the induced folding $\overline{F:} \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right) \longrightarrow \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$ has the following forms:

$$
\begin{aligned}
& \left.\left.\overline{F( } \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right)=\overline{F( } \pi_{1}\left(T_{1}^{1}\right)\right) * \pi_{1}\left(T_{2}^{1}\right) \\
& \left.\left.\overline{F( } \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right)=\pi_{1}\left(T_{1}^{1}\right) * \overline{F( } \pi_{1}\left(T_{2}^{1}\right)\right) \\
& \left.\left.\left.\overline{F( } \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right)=\overline{F( } \pi_{1}\left(T_{1}^{1}\right)\right) * \overline{F( } \pi_{1}\left(T_{2}^{1}\right)\right)
\end{aligned}
$$

Theorem 6 If $F_{i}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}, i=1,2$ are two types of foldings, where $F_{i}\left(T_{j}^{1}\right)=$ $T_{i}, j=1,2$. Then, there are induced foldings $\overline{F_{i}: \pi_{1}}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right) \longrightarrow \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$ such that $\overline{F_{i}}\left(\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right) \approx Z \times Z$.

Proof First, if $F_{1}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}$ is a folding such that $F_{1}\left(T_{1}^{1}\right)=T_{1}, F_{1}\left(T_{2}^{1}\right)=T_{1}$ as in Fig.6. Then we can express each element $g=a_{1} a_{2} \ldots a_{m}, m \geq 1$ of $\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$ in the following forms.

$$
\begin{aligned}
& \left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right)\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]^{k_{m-1}}\right)\left(\left[\alpha_{2}\right]^{n_{m}}\left[\beta_{2}\right]^{k_{m}}\right), \\
& \left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right)\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{2}\right]^{n_{m-1}},\left[\beta_{2}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right), \\
& \left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right)\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right)\left(\left[\alpha_{2}\right]^{n_{3}},\left[\beta_{2}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{2}\right]^{n_{m-1}},\left[\beta_{2}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right), \\
& \left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right)\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right)\left(\left[\alpha_{2}\right]^{n_{3}},\left[\beta_{2}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right),
\end{aligned}
$$

where $n_{1}, n_{2}, \cdots, n_{m}, k_{1}, k_{2}, \cdots, k_{m}$ are nonzero integers,

$$
\left(\left[\alpha_{i}\right]^{n_{1}},\left[\beta_{i}\right]^{k_{1}}\right) \in \pi_{1}\left(T_{1}^{1}\right),\left(\left[\alpha_{i}\right]^{n_{2}},\left[\beta_{i}\right]^{k_{2}}\right) \in \pi_{1}\left(T_{2}^{1}\right)
$$

Since $\overline{F_{1}}\left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)=\left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right), \overline{F_{1}}\left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right)=\left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)$, it follows that there is an induced folding $\overline{F_{i}}: \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right) \longrightarrow \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$ such that $\overline{F_{1}}\left(\pi_{1}\left(T_{1}^{1}\right) *\right.$ $\left.\pi_{1}\left(T_{2}^{1}\right)\right)=\pi_{1}\left(T_{1}^{1}\right)$, and so $\overline{F_{1}}\left(\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right) \approx Z \times Z$. Similarly, if $F_{2}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow$
$T_{1}^{1} \vee T_{2}^{1}$ is folding, such that $F_{2}\left(T_{1}^{1}\right)=T_{1}, F_{2}\left(T_{2}^{1}\right)=T_{1}$, then there is an induced folding $\overline{F_{2}}\left(\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right)=\pi_{1}\left(T_{1}^{1}\right)$ such that $\overline{F_{2}}\left(\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right) \approx Z \times Z$.


Fig. 6

Theorem 7 If $F_{i}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}, i=1,2$ are two types of foldings, where $F_{i}\left(T_{j}^{1}\right)=$ $T_{i}, j=1,2$. Then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right) \approx Z \times Z$.

Proof If $F_{i}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}, i=1,2$ are two types of foldings, where $F_{i}\left(T_{j}^{1}\right)=$ $T_{i}, j=1,2$, then $\lim _{n \rightarrow \infty} F_{i_{n}}\left(T_{1}^{1} \vee T_{2}^{1}\right)=T_{i}^{1}$ as in Fig.7. Thus $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)=\pi_{1}\left(T_{i}^{1}\right)$ ,since $\pi_{1}\left(T_{i}^{1}\right) \approx Z \times Z$ we have $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right) \approx Z \times Z$.


Fig. 7

Corollary 2 If $F_{i}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}, i=1,2$ are two types of foldings, where $F_{i}\left(T_{j}^{1}\right)=$ $T_{i}, j=1,2$. Then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$ is a free Abelian group of rank $2 n$.

Proof Since $F_{i}\left(T_{j}^{1}\right)=T_{i}, j=1,2$ we have the following chain $T_{1}^{1} \vee T_{2}^{1} \xrightarrow{F_{i_{1}}} T_{i}^{n} \xrightarrow{F_{i_{2}}}$ $T_{i}^{n} \ldots . \xrightarrow{\lim _{n \rightarrow \infty} F_{i n}} T_{i}^{n}$ Since $\pi_{1}\left(T_{i}^{n}\right)=\underbrace{\pi_{1}\left(T_{i} \times T_{i} \times \ldots . \times T_{i}\right)}_{n \text {-terms }}$, , it follows that $\pi_{1}\left(T_{i}^{n}\right) \approx \underbrace{Z \times Z \times \ldots \times Z}_{2 n-\text { terms }}$. Hence,$\pi_{1}\left(\lim _{n \rightarrow \infty} F_{i_{n}}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$ is a free Abelian of rank $2 n$.

Theorem 8 If $F: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}$ is a folding by cut such that $F_{1}\left(T_{1}^{1}\right) \neq T_{1}, F_{1}\left(T_{2}^{1}\right) \neq T_{1}$ .Then there is induced folding $\overline{F:} \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right) \longrightarrow \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$ such that $\bar{F}\left(\pi_{1}\left(T_{1}^{1}\right) *\right.$ $\left.\pi_{1}\left(T_{2}^{1}\right)\right) \approx Z * Z$.

Proof Let $F: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}$ is a folding such that $F_{1}\left(T_{1}^{1}\right) \neq T_{1}, F_{1}\left(T_{2}^{1}\right) \neq T_{1}$ as in Fig.8. Then, we can express each element $g=a_{1} a_{2} \cdots a_{m}, m \geq 1$ of $\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$ in the
following forms

$$
\begin{aligned}
& \left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right)\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]^{k_{m-1}}\right)\left(\left[\alpha_{2}\right]^{n_{m}}\left[\beta_{2}\right]^{k_{m}}\right), \\
& \left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right)\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{2}\right]^{n_{m-1}},\left[\beta_{2}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right), \\
& \left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right)\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right)\left(\left[\alpha_{2}\right]^{n_{3}},\left[\beta_{2}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{2}\right]^{n_{m-1}},\left[\beta_{2}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right), \\
& \left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right)\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right)\left(\left[\alpha_{2}\right]^{n_{3}},\left[\beta_{2}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right),
\end{aligned}
$$

where $n_{1}, n_{2}, \cdots, n_{m}, k_{1}, k_{2}, \cdots, k_{m}$ are nonzero integers and

$$
\left(\left[\alpha_{i}\right]^{n_{1}},\left[\beta_{i}\right]^{k_{1}}\right) \in \pi_{1}\left(T_{1}^{1}\right),\left(\left[\alpha_{i}\right]^{n_{2}},\left[\beta_{i}\right]^{k_{2}}\right) \in \pi_{1}\left(T_{2}^{1}\right) .
$$

Then, the induced folding of the element $g$ is

$$
\begin{aligned}
\bar{F}(g)= & \bar{F}\left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right) \bar{F}\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right) \bar{F}\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]{ }^{k_{3}}\right) \\
& \cdots \bar{F}\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]^{k_{m-1}}\right) \bar{F}\left(\left[\alpha_{2}\right]^{n_{m}}\left[\beta_{2}\right]^{k_{m}}\right) \\
= & \left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right)\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]^{k_{m-1}}\right)\left(\left[\alpha_{2}\right]^{n_{m}}\left[\beta_{2}\right]^{k_{m}}\right), \\
& \bar{F}\left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right) \bar{F}\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right) \bar{F}\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]^{k_{3}}\right) \\
& \cdots \bar{F}\left(\left[\alpha_{2}\right]^{n_{m}-1},\left[\beta_{2}\right]^{k_{m-1}}\right) \bar{F}\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right) \\
= & \left(\left[\alpha_{1}\right]^{n_{1}},\left[\beta_{1}\right]^{k_{1}}\right)\left(\left[\alpha_{2}\right]^{n_{2}},\left[\beta_{2}\right]^{k_{2}}\right)\left(\left[\alpha_{1}\right]^{n_{3}},\left[\beta_{1}\right]^{k_{3}}\right) \cdots\left(\left[\alpha_{2}\right]^{n_{m-1}},\left[\beta_{2}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right), \\
& \bar{F}\left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right) \bar{F}\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right) \bar{F}\left(\left[\alpha_{2}\right]^{n_{3}},\left[\beta_{2}\right]^{k_{3}}\right) \\
& \cdots \bar{F}\left(\left[\alpha_{2}\right]^{n_{m-1}},\left[\beta_{2}\right]_{m-1}^{k_{m-1}}\right) \bar{F}\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right) \\
= & \left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right)\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right)\left(\left[\alpha_{2}\right]^{n_{3}},\left[\beta_{2}\right]_{3}^{k_{3}}\right) \cdots\left(\left[\alpha_{2}\right]^{n_{m-1}},\left[\beta_{2}\right]_{m-1}^{m_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]_{m}\right), \\
& \bar{F}\left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right) \bar{F}\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right) \bar{F}\left(\left[\alpha_{2}\right]_{3},\left[\beta_{2}\right]^{k_{3}}\right) \\
& \cdots \bar{F}\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]_{m_{m-1}}\right) \bar{F}\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]_{m}\right) \\
= & \left(\left[\alpha_{2}\right]^{n_{1}},\left[\beta_{2}\right]^{k_{1}}\right)\left(\left[\alpha_{1}\right]^{n_{2}},\left[\beta_{1}\right]^{k_{2}}\right)\left(\left[\alpha_{2}\right]^{n_{3}},\left[\beta_{2}\right]_{3}^{k_{3}}\right) \cdots\left(\left[\alpha_{1}\right]^{n_{m-1}},\left[\beta_{1}\right]^{k_{m-1}}\right)\left(\left[\alpha_{1}\right]^{n_{m}}\left[\beta_{1}\right]^{k_{m}}\right) .
\end{aligned}
$$

It follows from $\left[\widehat{\beta}_{1}\right],\left[\widehat{\beta}_{2}\right] \rightarrow 0$ ( identity element) that there is an induced folding such that $\overline{F:} \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right) \longrightarrow \pi_{1}\left(S_{1}^{1}\right) * \pi_{1}\left(S_{2}^{1}\right)$. Therefore, $\bar{F}\left(\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right) \approx Z * Z$.


Fig. 8
Corollary 3 If $F_{i}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}, i=1,2$ are two types of foldings such that $F_{i}\left(T_{j}^{1}\right)=$ $T_{i}^{1}, F_{j}\left(T_{i}^{1}\right) \neq T_{i}^{1}, i, j=1,2, i \neq j$. Then there are induced foldings $\overline{F_{i}}: \pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right) \longrightarrow$ $\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)$. such that $\overline{F_{i}}\left(\pi_{1}\left(T_{1}^{1}\right) * \pi_{1}\left(T_{2}^{1}\right)\right) \approx(Z \times Z) * Z$.

Theorem 9 If $F: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}$ are a folding by cut such that $F\left(T_{i}^{1}\right) \neq T_{i}^{1}$, for $i=1,2$ . Then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$, is a free group of rank $\leq 2$ or identity group.

Proof Consider, $F\left(T_{i}^{1}\right) \neq T_{i}^{1}$, for $i=1,2$, then we have the following: $\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)=$ $S_{1}^{1} \vee S_{2}^{1}$ as in Fig.9(a) then,$\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right) \approx \pi_{1}\left(S_{1}^{1}\right) \vee \pi_{1}\left(S_{2}^{1}\right)$, and so $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$
$\approx Z * Z$. Hence, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$ is a free group of rank 2.Also, If $\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)$ as in Fig.9(b), then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)=0$. Moreover, if $\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)$ as in Fig.9(c), then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right) \approx \pi_{1}\left(S_{1}^{1}\right) \approx Z$.Therefore, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$ is a free group of rank $\leq 2$ or identity group.


(a)
or $\cdots \xrightarrow{\lim _{n \rightarrow \infty} F_{n}}$

(b)

(c)

Fig. 9

(a)
or


Fig. 10

Theorem 10 If $F_{i}: T_{1}^{1} \vee T_{2}^{1} \longrightarrow T_{1}^{1} \vee T_{2}^{1}, i=1,2$ are two types of foldings such that
$F_{i}\left(T_{i}^{1}\right)=T_{i}^{1}, F_{j}\left(T_{i}^{1}\right) \neq T_{i}^{1}, i, j=1,2, i \neq j$. Then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$ is either isomorphic $(Z \times Z) * Z$ to or $(Z \times Z)$.

Proof Since $F_{i}\left(T_{i}^{1}\right)=T_{i}^{1}, F_{j}\left(T_{i}^{1}\right) \neq T_{i}^{1}, i, j=1,2, i \neq j$, we have the following:
If $\lim _{n \rightarrow \infty} F_{i_{n}}\left(T_{1}^{1} \vee T_{2}^{1}\right)=T_{i}^{1} \vee S_{i}^{1}$ as in Fig.10(a), then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)=\pi_{1}\left(T_{i}^{1} \vee S_{i}^{1}\right) \approx$ $(Z \times Z) * Z$. Also, if $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)=\pi_{1}\left(T_{i}^{1}\right)$ as in Fig.10(b) then $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee\right.\right.$ $\left.T_{2}^{1}\right) \pi_{1}\left(T_{i}^{1}\right) \approx Z \times Z$. Hence, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{1} \vee T_{2}^{1}\right)\right)$ is either isomorphic to $(Z \times Z) * Z$ or $(Z \times Z)$.

Theorem 11 If $F: T_{1}^{n} \vee T_{2}^{n} \longrightarrow T_{1}^{n} \vee T_{2}^{n}$ is a folding such that $F\left(T_{1}^{n}\right)=T_{1}^{n}$ and $\quad F\left(T_{2}^{n}\right) \neq$ $T_{2}^{n}$ where $F\left(T_{2}^{n}\right)=\underbrace{F\left(T_{2}^{1}\right) \times F\left(T_{2}^{1}\right) \times \ldots \times F\left(T_{2}^{1}\right)}_{n \text {-terms }}, F\left(T_{2}^{1}\right) \neq T_{2}^{1}$ is a folding by cut. Then, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{n} \vee T_{2}^{n}\right)\right)$ is isomorphic to $\underbrace{(Z \times Z \times \ldots \times Z}_{2 n-\text { terms }}) * \underbrace{Z \times Z \times \ldots \times Z}_{n-\text { terms }}$.

Proof Since $F\left(T_{1}^{n}\right)=T_{1}^{n}, F\left(T_{2}^{n}\right) \neq T_{2}^{n}$ we have the following chain:

$$
\begin{aligned}
& T_{1}^{n} \vee T_{2}^{n} \xrightarrow{F} T_{1}^{n} \vee \underbrace{F\left(S_{1}^{1}\right) \times S_{2}^{1} \times F\left(S_{1}^{1}\right) \times S_{2}^{1} \times \cdots \times F\left(S_{1}^{1}\right) \times S_{2}^{1}}_{2 n-\text { terms }} \stackrel{F}{\rightarrow}, \\
& T_{1}^{n} \vee T_{2}^{n} \xrightarrow{F} T_{1}^{n} \vee \underbrace{F\left(S_{1}^{1}\right) \times S_{2}^{1} \times F\left(S_{1}^{1}\right) \times S_{2}^{1} \times \cdots \times F\left(S_{1}^{1}\right) \times S_{2}^{1}}_{2 n-\text { terms }} \stackrel{F}{\rightarrow}, \\
& T_{1}^{n} \vee \underbrace{F\left(F\left(S_{1}^{1}\right)\right) \times S_{2}^{1} \times F\left(F\left(S_{1}^{1}\right)\right) \times S_{2}^{1} \times \cdots \times F\left(F\left(S_{1}^{1}\right)\right) \times S_{2}^{1}}_{2 n-\text { terms }}{ }_{\underline{n \rightarrow \infty}}^{\lim _{n \rightarrow \infty} F_{n},} \\
& T_{1}^{n} \vee \underbrace{\left(S_{2}^{1} \times S_{2}^{1} \times \cdots \times S_{2}^{1}\right)}_{n-\text { terms }} .
\end{aligned}
$$

Hence, $\pi_{1}\left(\lim _{n \rightarrow \infty} F_{n}\left(T_{1}^{n} \vee T_{2}^{n}\right)\right)$ is isomorphic to $(\underbrace{Z \times Z \times \cdots \times Z}_{2 n-\text { terms }}) * \underbrace{Z \times Z \times \cdots \times Z}_{n-\text { terms }}$.
Theorem 12 Let $F: M \rightarrow M$ is a folding by cut or with singularity, and $M$ is a manifold homeomorphic to $S^{1}$ or $T^{1}$. Then, there are unfoldings unf $: F(M) \subset M \rightarrow M$ such that $\pi_{1}\left(\lim _{n \rightarrow \infty} u n f_{n}(F(M))\right.$ is isomorphic to $Z$ or $Z \times Z$.

Proof We have two cases following.
Case 1. Let $M$ be a manifold homeomorphic to $S^{1}$, if $F: S^{1} \rightarrow S^{1}$ is a folding by cut.


## Fig. 11

Then, we can define a sequence of unfoldings

$$
\begin{aligned}
& u n f_{1}: F\left(S^{1}\right) \rightarrow M_{1}, F\left(S^{1}\right) \neq S^{1}, M_{1} \subseteq S^{1}, u n f_{2}: M_{1} \rightarrow M_{2}, \ldots, u n f_{n}: M_{1} \rightarrow M_{2} \\
& \lim _{n \rightarrow \infty} \operatorname{unf} f_{n}(F(M))=S^{1} \text { as in Fig.11. Thus } \pi_{1}\left(\lim _{n \rightarrow \infty} \operatorname{unf}_{n}(F(M)) \approx Z .\right.
\end{aligned}
$$

Case 2. Let $M$ be a manifold homeomorphic to $T^{1}$, if $F: T^{1} \rightarrow T^{1}$ is a folding such that $F\left(S_{1}^{1}\right)=S_{1}^{1}$ and $F\left(S_{2}^{1}\right) \neq S_{2}^{1}$. So we can define a sequence of unfoldings following.
$u n f_{1}: F\left(T^{1}\right) \rightarrow M_{1}, u n f_{2}: M_{1} \rightarrow M_{2}, \cdots, u n f_{n}: M_{1} \rightarrow M_{2}$,
$\lim _{n \rightarrow \infty} \operatorname{unf}_{n}(F(M))=T^{1}$ as in Fig.12. Thus $\pi_{1}\left(\lim _{n \rightarrow \infty} \operatorname{unf}_{n}(F(M)) \approx Z \times Z\right.$.


Fig. 12
Therefore, $\pi_{1}\left(\lim _{n \rightarrow \infty} u n f_{n}(F(M))\right.$ is isomorphic to $Z$ or $Z \times Z$.
Corollary 4 Let $F: M \rightarrow M$ be a folding by cut or with singularity, $M$ is a manifold homeomorphic to $S^{n}$ or $T^{n}, n \geq 2$. Then there are unfoldings unf $: F(M) \subset M \rightarrow M$ such that $\pi_{1}\left(\lim _{n \rightarrow \infty} \operatorname{unf}_{n}(F(M))\right.$ is the identity group or a free Abelian group of rank $2 n$.

## References

[1] P.DI-Francesco, Folding and coloring problem in mathematics and physics, Bulletin of the American Mathematical Society, Vol. 37, No. 3 (2000), 251-307.
[2] A.E.El-Ahmady, The deformation retract and topological folding of Buchdahi space, Periodica Mathematica Hungarica, Vol. 28. No. 1(1994),19-30.
[3] A.E.El-Ahmady, Fuzzy Lobacherskian space and its folding, The Journal of Fuzzy Mathematics, Vol. 2, No.2, (2004),255-260.
[4] A.E.El-Ahmady: Fuzzy folding of fuzzy horocycle, Circolo Mathematico Palermo, Serie II, Tomo LIII (2004), 443-450.
[5] M.El-Ghoul, Unfolding of Riemannian manifolds, Commun. Fac. Sci. Univ Ankara, Series, A37 (1988), 1-4.
[6] M.El- Ghoul, The deformation retract of the complex projective space and its topological folding, Journal of Material Science, Vol. 30 (1995), 4145-4148.
[7] M.El-Ghoul, Fractional folding of a manifold, Chaos Solitons and Fractals, Vol. 12 (2001), 1019-1023.
[8] M.El-Ghoul, A.E.El-Ahmady, and H.Rafat, Folding-retraction of chaotic dynamical manifold and the VAK of vacuum fluctation, Chaos Solutions and Fractals, Vol. 20 (2004), 209-217.
[9] E.El-Kholy, Isometric and Topological Folding of Manifold, Ph. D. Thesis, University of Southampton, UK (1981).
[10] R.Frigerio, Hyperbolic manifold with geodesic boundary which are determine by their fundamental group, Topology and its Application, 45 (2004), 69-81.
[11] A.Hatcher, Algebraic Topology, The web address is: http: /www.math.coronell.edu/hatcher.
[12] W.S.Massey, Algebraic Topology: An Introduction, Harcourt Brace and world, New York (1967).
[13] O.Neto and P.C.Silva, The fundamental group of an algebraic link, C. R. Cad Sci. Paris, Ser. I, 340 (2005), 141-146.
[14] S.A.Robertson, Isometric folding of Riemannian manifolds, Proc. Roy. Soc. Edinburgh, 77 (1977), 275-289.

# Absolutely Harmonious Labeling of Graphs 

M.Seenivasan<br>(Sri Paramakalyani College, Alwarkurichi-627412, India)<br>A.Lourdusamy<br>(St.Xavier's College (Autonomous), Palayamkottai, India)<br>E-mail: msvasan_22@yahoo.com, lourdugnanam@hotmail.com


#### Abstract

Absolutely harmonious labeling $f$ is an injection from the vertex set of a graph $G$ with $q$ edges to the set $\{0,1,2, \ldots, q-1\}$, if each edge $u v$ is assigned $f(u)+f(v)$ then the resulting edge labels can be arranged as $a_{0}, a_{1}, a_{2}, \ldots, a_{q-1}$ where $a_{i}=q-i$ or $q+i, 0 \leq$ $i \leq q-1$. However, when $G$ is a tree one of the vertex labels may be assigned to exactly two vertices. A graph which admits absolutely harmonious labeling is called absolutely harmonious graph. In this paper, we obtain necessary conditions for a graph to be absolutely harmonious and study absolutely harmonious behavior of certain classes of graphs.


Key Words: Graph labeling, Smarandachely $k$-labeling, harmonious labeling, absolutely harmonious labeling.

## AMS(2010): O5C78

## §1. Introduction

A vertex labeling of a graph $G$ is an assignment $f$ of labels to the vertices of $G$ that induces a label for each edge $x y$ depending on the vertex labels. For an integer $k \geq 1$, a Smarandachely $k$-labeling of a graph $G$ is a bijective mapping $f: V \rightarrow\{1,2, \cdots, k|V(G)|+|E(G)|\}$ with an additional condition that $|f(u)-f(v)| \geq k$ for $\forall u v \in E$. particularly, if $k=1$, i.e., such a Smarandachely 1-labeling is the usually labeling of graph. Among them, labelings such as those of graceful labeling, harmonious labeling and mean labeling are some of the interesting vertex labelings found in the dynamic survey of graph labeling by Gallian [2]. Harmonious labeling is one of the fundamental labelings introduced by Graham and Sloane [3] in 1980 in connection with their study on error correcting code. Harmonious labeling $f$ is an injection from the vertex set of a graph $G$ with $q$ edges to the set $\{0,1,2, \ldots, q-1\}$, if each edge $u v$ is assigned $f(u)+f(v)(\bmod q)$ then the resulting edge labels are distinct. However, when $G$ is a tree one of the vertex labels may be assigned to exactly two vertices. Subsequently a few variations of harmonious labeling, namely, strongly c-harmonious labeling [1], sequential labeling [5], elegant labeling [1] and felicitous labeling [4] were introduced. The later three labelings were introduced to avoid such exceptions for the trees given in harmonious labeling. A strongly

[^4]1-harmonious graph is also known as strongly harmonious graph.
It is interesting to note that if a graph $G$ with $q$ edges is harmonious then the resulting edge labels can be arranged as $a_{0}, a_{1}, a_{2}, \cdots, a_{q-1}$ where $a_{i}=i$ or $q+i, 0 \leq i \leq q-1$. That is for each $i, 0 \leq i \leq q-1$ there is a distinct edge with label either $i$ or $q+i$. An another interesting and natural variation of edge label could be $q-i$ or $q+i$. This prompts to define a new variation of harmonious labeling called absolutely harmonious labeling.

Definition 1.1 An absolutely harmonious labeling $f$ is an injection from the vertex set of a graph $G$ with $q$ edges to the set $\{0,1,2, \ldots, q-1\}$, if each edge uv is assigned $f(u)+f(v)$ then the resulting edge labels can be arranged as $a_{0}, a_{1}, a_{2}, \ldots, a_{q-1}$ where $a_{i}=q-i$ or $q+i, 0 \leq i \leq q-1$. However, when $G$ is a tree one of the vertex labels may be assigned to exactly two vertices. A graph which admits absolutely harmonious labeling is called absolutely harmonious graph.

The result of Graham and Sloane [3] states that $C_{n}, n \cong 1(\bmod 4)$ is harmonious, but we show that $C_{n}, n \cong 1(\bmod 4)$ is not an absolutely harmonious graph. On the other hand, we show that $C_{4}$ is an absolutely harmonious graph, but it is not harmonious. We observe that a strongly harmonious graph is an absolutely harmonious graph.

To initiate the investigation on absolutely harmonious graphs, we obtain necessary conditions for a graph to be an absolutely harmonious graph and prove the following results:

1. Path $P_{n}, n \geq 3$, a class of banana trees, and $P_{n} \odot K_{m}^{c}$ are absolutely harmonious graphs.
2. Ladders, $C_{n} \odot K_{m}^{c}$, Triangular snakes, Quadrilateral snakes, and $m K_{4}$ are absolutely harmonious graphs.
3. Complete graph $K_{n}$ is absolutely harmonious if and only if $n=3$ or 4 .
4. Cycle $C_{n}, n \cong 1$ or $2(\bmod 4), C_{m} \times C_{n}$ where $m$ and $n$ are odd, $m K_{3}, m \geq 2$ are not absolutely harmonious graphs.

## §2. Necessary Conditions

Theorem 2.1 If $G$ is an absolutely harmonious graph, then there exists a partition $\left(V_{1}, V_{2}\right)$ of the vertex set $V(G)$, such that the number of edges connecting the vertices of $V_{1}$ to the vertices of $V_{2}$ is exactly $\left\lceil\frac{q}{2}\right\rceil$.

Proof If $G$ is an absolutely harmonious graph, then the vertices can be partitioned into two sets $V_{1}$ and $V_{2}$ having respectively even and odd vertex labels. Observe that among the $q$ edges $\frac{q}{2}$ edges or $\left\lceil\frac{q}{2}\right\rceil$ edges are labeled with odd numbers, according as $q$ is even or $q$ is odd. For an edge to have odd label, one end vertex must be odd labeled and the other end vertex must be even labeled. Thus, the number of edges connecting the vertices of $V_{1}$ to the vertices of $V_{2}$ is exactly $\left\lceil\frac{q}{2}\right\rceil$.

Remark 2.2 A simple and straight forward application of Theorem 2.1 identifies the non absolutely harmonious graphs. For example, complete graph $K_{n}$ has $\frac{n(n-1)}{2}$ edges. If we assign
$m$ vertices to the part $V_{1}$, there will be $m(n-m)$ edges connecting the vertices of $V_{1}$ to the vertices of $V_{2}$. If $K_{n}$ has an absolutely harmonious labeling, then there is a choice of $m$ for which $m(n-m)=\left\lceil\frac{n^{2}-n}{4}\right\rceil$. Such a choice of $m$ does not exist for $n=5,7,8.10, \ldots$.

A graph is called even graph if degree of each vertex is even.
Theorem 2.3 If an even graph $G$ is absolutely harmonious then $q \cong 0 \operatorname{or} 3(\bmod 4)$.
Proof Let $G$ be an even graph with $q \cong 1$ or $2(\bmod 4)$ and $d(v)$ denotes the degree of the vertex $v$ in $G$. Suppose $f$ be an absolutely harmonious labeling of $G$. Then the resulting edge labels can be arranged as $a_{0}, a_{1}, a_{2}, \ldots, a_{q-1}$ where $a_{i}=q-i$ or $q+i, 0 \leq i \leq q-1$. In other words, for each $i$, the edge label $a_{i}$ is $(q-i)+2 i b_{i}, 0 \leq i \leq q-1$ where $b_{i} \in\{0,1\}$. Evidently

$$
\sum_{v \in V(G)} d(v) f(v)-2 \sum_{k=0}^{q-1} k b_{k}=\binom{q+1}{2}
$$

As $d(v)$ is even for each $v$ and $q \cong 1$ or $2(\bmod 4)$,

$$
\sum_{v \in V(G)} d(v) f(v)-2 \sum_{k=0}^{q-1} k b_{k} \cong 0(\bmod 2)
$$

but $\binom{q+1}{2} \cong 1(\bmod 2)$. This contradiction proves the theorem.
Corollary $2.4 A$ cycle $C_{n}$ is not an absolutely harmonious graph if $n \cong 1 \operatorname{or} 2(\bmod 4)$.
Corollary 2.5 $A$ grid $C_{m} \times C_{n}$ is not an absolutely harmonious graph if $m$ and $n$ are odd.
Theorem 2.6 If $f$ is an absolutely harmonious labeling of the cycle $C_{n}$, then edges of $C_{n}$ can be partitioned into two sub sets $E_{1}, E_{2}$ such that

$$
\sum_{u v \in E_{1}}|f(u)+f(v)-n|=\frac{n(n+1)}{4} \text { and } \sum_{u v \in E_{2}}|f(u)+f(v)-n|=\frac{n(n-3)}{4}
$$

Proof Let $v_{1} v_{2} v_{3} \ldots v_{n} v_{1}$ be the cycle $C_{n}$, where $e_{i}=v_{i-1} v_{i}, 2 \leq i \leq n$ and $e_{1}=v_{n} v_{1}$. Define $E_{1}=\{u v \in E / f(u)+f(v)-n$ is non negative $\}$ and $E_{2}=\{u v \in E / f(u)+f(v)-$ $n$ is negative $\}$. Since $f$ is an absolutely harmonious labeling of the cycle $C_{n}$,

$$
\sum_{u v \in E}|f(u)+f(v)-n|=\frac{n(n-1)}{2}
$$

In other words,

$$
\begin{equation*}
\sum_{u v \in E_{1}}|f(u)+f(v)-n|+\sum_{u v \in E_{2}}|f(u)+f(v)-n|=\frac{n(n-1)}{2} \tag{1}
\end{equation*}
$$

Since $\sum_{u v \in E}(f(u)+f(v)-n)=-n$, we have

$$
\begin{equation*}
\sum_{u v \in E_{1}}|f(u)+f(v)-n|-\sum_{u v \in E_{2}}|f(u)+f(v)-n|=-n \tag{2}
\end{equation*}
$$

Solving equations (1) and (2), we get the desired result.
Remark 2.7 If $n \cong 1$ or $2(\bmod 4)$ then both $\frac{n(n+1)}{4}$ and $\frac{n(n-3)}{4}$ cannot be integers. Thus the cycle $C_{n}$ is not an absolutely harmonious graph if $n \cong 1$ or ${ }_{2}^{4}(\bmod 4)$.

Remark 2.8 Observe that the conditions stated in Theorem 2.1, Theorem 2.3, and Theorem 2.6 are necessary but not sufficient. Note that $C_{8}$ satisfies all the conditions stated in Theorems $2.1,2.3$, and 2.6 but it is not an absolutely harmonious graph. For, checking each of the $\frac{8!}{2}$ possibilities reveals the desired result about $C_{8}$.

## §3. Absolutely Harmonious Graphs

Theorem 3.1 The path $P_{n+1}$, where $n \geq 2$ is an absolutely harmonious graph.
Proof Let $P_{n+1}: v_{1} v_{2} \ldots v_{n+1}$ be a path, $r=\left\lceil\frac{n}{2}\right\rceil, s=\left\{\begin{array}{ll}\left\lceil\frac{r}{2}\right\rceil+1 & \text { if } n \cong 0(\bmod 4) \\ \left\lceil\frac{r}{2}\right\rceil & \text { otherwise }\end{array}\right.$, $t=\left\{\begin{array}{ll}s-1 & \text { if } n \cong 0 \text { or } 1(\bmod 4) \\ s & \text { if } n \cong 2 \operatorname{or} 3(\bmod 4)\end{array}, T_{1}=n, T_{2}=\left\{\begin{array}{ll}2 t+2 & \text { if } n \cong 0 \operatorname{or} 1(\bmod 4) \\ 2 t+1 & \text { if } n \cong 2 \operatorname{or} 3(\bmod 4)\end{array}\right.\right.$ and $T_{3}=$ $\left\{\begin{array}{ll}-1 & \text { if } n \cong 0 \text { or } 1(\bmod 4) \\ -2 & \text { if } n \cong 2 \text { or } 3(\bmod 4)\end{array}\right.$.

Then $r+s+t=n+1$. Define $f: V\left(P_{n+1}\right) \rightarrow\{0,1,2,3, \cdots, n-1\}$ by:
$f\left(v_{i}\right)=T_{1}-i$ if $1 \leq i \leq r, f\left(v_{r+i}\right)=T_{2}-2 i$ if $1 \leq i \leq s$ and $f\left(v_{r+s+i}\right)=T_{3}+2 i$ if $1 \leq i \leq t$.

Evidently $f$ is an absolutely harmonious labeling of $P_{n+1}$. For example, an absolutely harmonious labeling of $P_{12}$ is shown in Fig.3.1.


Fig. 3.1
The tree obtained by joining a new vertex $v$ to one pendant vertex of each of the $k$ disjoint stars $K_{1, n_{1}}, K_{1, n_{2}}, K_{1, n_{3}}, \ldots, K_{1, n_{k}}$ is called a banana tree. The class of all such trees is denoted by $B T\left(n_{1}, n_{2}, n_{3}, \ldots, n_{k}\right)$.

Theorem 3.2 The banana tree $B T(n, n, n, \ldots, n)$ is absolutely harmonious.


Fig. 3.2
Proof Let $V\left(B T(n, n, n, \cdots, n)=\{v\} \cup\left\{v_{j}, v_{j r}: 1 \leq j \leq k\right.\right.$ and $\left.1 \leq r \leq n\right\}$ where $d\left(v_{j}\right)=$ $n$ and $E\left(B T(n, n, n, \ldots, n)=\left\{v v_{j n}: 1 \leq j \leq k\right\} \cup\left\{v_{j} v_{j r}: 1 \leq j \leq k, 1 \leq r \leq n\right\}\right.$. Clearly $B T(n, n$, $\cdots, n)$ has order $(n+1) k+1$ and size $(n+1) k$. Define

$$
f: V(B T(n, n, \cdots, n) \rightarrow\{1,2,3, \ldots,(n+1) k-1\}
$$

as follows:

$$
f(v)=1, f\left(v_{j}\right)=(n+1)(j-1): 1 \leq j \leq k, f\left(v_{j r}\right)=(n+1)(j-1)+r: 1 \leq r \leq n
$$

It can be easily verified that $f$ is an absolutely harmonious labeling of $B T(n, n, n, \ldots, n)$. For example an absolutely harmonious labeling of $B T(4,4,4,4)$ is shown in Fig.3.2.

The corona $G_{1} \odot G_{2}$ of two graphs $G_{1}\left(p_{1}, q_{1}\right)$ and $G_{2}\left(p_{2}, q_{2}\right)$ is defined as the graph obtained by taking one copy of $G_{1}$ and $p_{1}$ copies of $G_{2}$ and then joining the $i^{t h}$ vertex of $G_{1}$ to all the vertices in the $i^{t h}$ copy of $G_{2}$.

Theorem 3.3 The corona $P_{n} \odot K_{m}^{C}$ is absolutely harmonious.
Proof Let $V\left(P_{n} \odot K_{m}^{C}\right)=\left\{u_{i}: 1 \leq i \leq n\right\} \cup\left\{u_{i j}: 1 \leq i \leq n, 1 \leq j \leq m\right\}$ and $E\left(P_{n} \odot K_{m}^{C}\right)=$ $\left\{u_{i} u_{i+1}: 1 \leq i \leq n-1\right\} \cup\left\{u_{i} u_{i j}: 1 \leq i \leq n, 1 \leq j \leq m\right\}$. We observe that $P_{n} \odot K_{m}^{C}$ has order $(m+1) n$ and size $(m+1) n-1$. Define $f: V\left(P_{n} \odot K_{m}^{C}\right) \longrightarrow\{0,1,2, \ldots, m n+n-2\}$ as follows:

$$
f\left(u_{i}\right)=\left\{\begin{array}{ll}
0 & \text { if } i=1, \\
(m+1)(i-1) & \text { if } i=\left\lceil\frac{n}{2}\right\rceil \\
(m+1)(i-1)-1 & \text { otherwise }
\end{array} \quad \quad \quad\left(u_{i m}\right)= \begin{cases}(m+1) i & \text { if } 1 \leq i \leq\left\lceil\frac{n}{2}\right\rceil-2 \\
(m+1) i-1 & \text { if } i=\left\lceil\frac{n}{2}\right\rceil-1 \\
(m+1) i-2 & \left\lceil\frac{n}{2}\right\rceil \leq i \leq n\end{cases}\right.
$$

and for $1 \leq j \leq m-1$,

$$
f\left(u_{i j}\right)= \begin{cases}(m+1)(i-1)+j & \text { if } 1 \leq i \leq\left\lceil\frac{n}{2}\right\rceil-1, \\ (m+1)(i-1)+j-1 & \text { if }\left\lceil\frac{n}{2}\right\rceil \leq i \leq n .\end{cases}
$$

It can be easily verified that $f$ is an absolutely harmonious labeling of $P_{n} \odot K_{m}^{C}$. For example an absolutely harmonious labeling of $P_{5} \odot K_{3}^{C}$ is shown in Fig. 3.3.


Fig. 3.3
Theorem 3.4 The corona $C_{n} \odot K_{m}^{C}$ is absolutely harmonious.
Proof Let $V\left(C_{n} \odot K_{m}^{C}\right)=\left\{u_{i}: 1 \leq i \leq n\right\} \cup\left\{u_{i j}: 1 \leq i \leq n, 1 \leq j \leq m\right\}$ and $E\left(C_{n} \odot K_{m}^{C}\right)=$ $\left\{u_{i} u_{i+1}: 1 \leq i \leq n-1\right\} \cup\left\{u_{n} u_{1}\right\} \cup\left\{u_{i} u_{i j}: 1 \leq i \leq n, 1 \leq j \leq m\right\}$. We observe that $C_{n} \odot K_{m}^{C}$ has order $(m+1) n$ and size $(m+1) n$. Define $f: V\left(C_{n} \odot K_{m}^{C}\right) \longrightarrow\{0,1,2, \ldots, m n+n-1\}$ as follows:

$$
f\left(u_{i}\right)= \begin{cases}0 & \text { if } i=1 \\
(m+1)(i-1)-1 & \text { if } 2 \leq i \leq \frac{n-1}{2},, \quad f\left(u_{i m}\right)=\left\{\begin{array}{ll}
(m+1) i & \text { if } 1 \leq i \leq \frac{n-3}{2} \\
(m+1)(i-1) & \text { otherwise }
\end{array}, \quad\right. \text { otherwise } \\
(m+1) i-1 & \end{cases}
$$

and for $1 \leq j \leq m-1$

$$
f\left(u_{i j}\right)= \begin{cases}(m+1)(i-1)+j & \text { if } 1 \leq i \leq\left\lceil\frac{n}{2}\right\rceil-1 \\ (m+1)(i-1)+j-1 & \text { if }\left\lceil\frac{n}{2}\right\rceil \leq i \leq n\end{cases}
$$

It can be easily verified that $f$ is an absolutely harmonious labeling of $C_{n} \odot K_{m}^{C}$. For example an absolutely harmonious labeling of $C_{5} \odot K_{3}^{C}$ is shown in Figure 3.4.


Fig. 3.4
Theorem 3.5 The ladder $P_{n} \times P_{2}$, where $n \geq 2$ is an absolutely harmonious graph.
Proof Let $V\left(P_{n} \times P_{2}\right)=\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{n}\right\} \cup\left\{v_{1}, v_{2}, v_{3}, \ldots, v_{n}\right\}$ and $E\left(P_{n} \times P_{2}\right)=\left\{u_{i} u_{i+1}\right.$ : $1 \leq i \leq n-1\} \cup\left\{v_{i} v_{i+1}: 1 \leq i \leq n-1\right\} \cup\left\{u_{i} v_{i}: 1 \leq i \leq n\right\}$. We note that $P_{n} \times P_{2}$ has order
$2 n$ and size $3 n-2$.
Case 1. $n \equiv 0(\bmod 4)$.
Define $f: V\left(P_{n} \times P_{2}\right) \longrightarrow\{0,1,2, \ldots, 3 n-3\}$ by

$$
\begin{gathered}
f\left(u_{i}\right)= \begin{cases}3 i-2 & \text { if } i \text { is odd, } \\
3 i-2 & \text { if } i \text { is even and } 2 \leq i \leq \frac{n-4}{2}, \\
3 i-1 & \text { if } i \text { is even and } i=\frac{n}{2}, \\
3 i-3 & \text { if } i \text { is even and } \frac{n+4}{2} \leq i \leq n,\end{cases} \\
f\left(v_{1}\right)=0, \quad f\left(v_{\frac{n+2}{2}}\right)=\frac{3 n-6}{2}, \quad f\left(v_{i+1}\right)=f\left(u_{i}\right)+1 \text { if } 1 \leq i \leq n-1 \text { and } i \neq \frac{n}{2} .
\end{gathered}
$$

Case 2. $n \equiv 1(\bmod 4)$.
Define $f: V\left(P_{n} \times P_{2}\right) \longrightarrow\{0,1,2, \ldots, 3 n-3\}$ by

$$
\begin{gathered}
f\left(u_{i}\right)= \begin{cases}3 i-2 & \text { if } i \text { is odd and } 1 \leq i \leq \frac{n-3}{2}, \\
3 i-1 & \text { if } i=\frac{n+1}{2}, \\
3 i-3 & \text { if } i \text { is odd and } \frac{n+5}{2} \leq i \leq n, \\
3 i-2 & \text { if } i \text { is even, }\end{cases} \\
f\left(v_{1}\right)=0, \quad f\left(v_{\frac{n+3}{2}}\right)=\frac{3 n-3}{2}, \quad f\left(v_{i+1}\right)=f\left(u_{i}\right)+1 \text { if } 1 \leq i \leq n-1 \text { and } i \neq \frac{n+1}{2} .
\end{gathered}
$$

Case 3. $n \equiv 2(\bmod 4)$.
Define $f: V\left(P_{n} \times P_{2}\right) \longrightarrow\{0,1,2, \ldots, 3 n-3\}$ by

$$
\begin{aligned}
& f\left(u_{i}\right)= \begin{cases}3 i-2 & \text { if } i \text { is odd } \\
3 i-2 & \text { if } i \text { is even and } 2 \leq i \leq \frac{n-2}{2} \\
3 i-3 & \text { if } i \text { is even and } \frac{n+2}{2} \leq i \leq n\end{cases} \\
& f\left(v_{1}\right)=0, \quad f\left(v_{i+1}\right)=f\left(u_{i}\right)+1 \text { if } 1 \leq i \leq n-1
\end{aligned}
$$

Case 4. $n \equiv 3(\bmod 4)$.
Define $f: V\left(P_{n} \times P_{2}\right) \longrightarrow\{0,1,2, \ldots, 3 n-3\}$ by

$$
\begin{aligned}
& f\left(u_{i}\right)= \begin{cases}3 i-2 & \text { if } i \text { is odd and } 1 \leq i \leq \frac{n-1}{2} \\
3 i-3 & \text { if } i \text { is odd and } \frac{n+3}{2} \leq i \leq n \\
3 i-2 & \text { if } i \text { is even }\end{cases} \\
& f\left(v_{1}\right)=0, \quad f\left(v_{i+1}\right)=f\left(u_{i}\right)+1 \text { if } 1 \leq i \leq n-1 .
\end{aligned}
$$

In all four cases, it can be easily verified that $f$ is an absolutely harmonious labeling of $P_{n} \times P_{2}$. For example, an absolutely harmonious labeling of $P_{9} \times P_{2}$ is shown in Fig.3.5.


## Fig. 3.5

A $K_{n}$-snake has been defined as a connected graph in which all blocks are isomorphic to $K_{n}$ and the block-cut point graph is a path. A $K_{3}$-snake is called triangular snake.

Theorem 3.6 A triangular snake with $n$ blocks is absolutely harmonious if and only if $n \cong$ 0 or $1(\bmod 4)$.

Proof The necessity follows from Theorem 2.3.Let $G_{n}$ be a triangular snake with $n$ blocks on $p$ vertices and $q$ edges. Then $p=2 n-1$ and $q=3 n$. Let $V\left(G_{n}\right)=\left\{u_{i}: 1 \leq i \leq n+1\right\} \cup$ $\left\{v_{i}: 1 \leq i \leq n\right\}$ and $E\left(G_{n}\right)=\left\{u_{i} u_{i+1}, u_{i} v_{i}, u_{i+1} v_{i}: 1 \leq i \leq n\right\}$.

Case 1. $\quad n \equiv 0(\bmod 4)$.
Let $m=\frac{n}{4}$. Define $f: V\left(G_{n}\right) \longrightarrow\{0,1,2, \ldots, 3 n-1\}$ as follows:

$$
\begin{aligned}
& f\left(u_{i}\right)= \begin{cases}0 & \text { if } i=1 \\
2 i-2 & \text { if } 2 \leq i \leq 3 m \text { and } i \equiv 0 \operatorname{or} 2(\bmod 3) \\
2 i-1 & \text { if } 2 \leq i \leq 3 m \text { and } i \equiv 1(\bmod 3) \\
6 i-3 n-7 & \text { otherwise }\end{cases} \\
& f\left(v_{i}\right)= \begin{cases}1 & \text { if } i=1, \\
2 i-1 & \text { if } 2 \leq i \leq 3 m-1 \text { and } i \equiv \operatorname{or} 2(\bmod 3) \\
2 i-2 & \text { if } 2 \leq i \leq 3 m-1 \text { and } i \equiv 1(\bmod 3) \\
6 m+1 & \text { if } i=3 m \\
6 i-3 n-3 & \text { otherwise }\end{cases}
\end{aligned}
$$

Case 2. $\quad n \equiv 1(\bmod 4)$.
Let $m=\frac{n-1}{4}$. Define $f: V\left(G_{n}\right) \longrightarrow\{0,1,2, \ldots, 3 n-1\}$ as follows:

$$
\begin{aligned}
& f\left(u_{i}\right)= \begin{cases}0 & \text { if } i=1, \\
2 i-2 & \text { if } 2 \leq i \leq 3 m+2 \text { and } i \equiv 0 \text { or } 2(\bmod 3), \\
2 i-1 & \text { if } 2 \leq i \leq 3 m+2 \text { and } i \equiv 1(\bmod 3), \\
6 i-3 n-7 & \text { otherwise },\end{cases} \\
& f\left(v_{i}\right)= \begin{cases}1 & \text { if } i=1, \\
2 i-1 & \text { if } 2 \leq i \leq 3 m+1 \text { and } i \equiv 0 \operatorname{or} 2(\bmod 3) \\
2 i-2 & \text { if } 2 \leq i \leq 3 m+1 \text { and } i \equiv 1(\bmod 3) \\
6 i-3 n-3 & \text { otherwise. }\end{cases}
\end{aligned}
$$

In both cases, it can be easily verified that $f$ is an absolutely harmonious labeling of the triangular snake $G_{n}$. For example, an absolutely harmonious labeling of a triangular snake with five blocks is shown in Fig.3.6.


Fig.3.6

Theorem 3.7 $K_{4}$-snakes are absolutely harmonious.

Proof Let $G_{n}$ be a $K_{4}$-snake with $n$ blocks on $p$ vertices and $q$ edges. Then $p=3 n+1$ and $q=6 n$. Let $V\left(G_{n}\right)=\left\{u_{i}, v_{i}, w_{i}: 1 \leq i \leq n\right\} \cup\left\{v_{n+1}\right\}$ and $E\left(G_{n}\right)=\left\{u_{i} v_{i}, u_{i} w_{i}, v_{i} w_{i}: 1 \leq i \leq n\right\} \cup$ $\left\{u_{i} v_{i+1}, v_{i} v_{i+1}, w_{i} v_{i+1}: 1 \leq i \leq n\right\}$ Define $f: V\left(G_{n}\right) \longrightarrow\{0,1,2, \ldots, 6 n-1\}$ as follows:

$$
f\left(u_{i}\right)=3 i-3, f\left(v_{i}\right)=3 i-2, f\left(w_{i}\right)=3 i-1
$$

where $1 \leq i \leq n$, and $f\left(v_{n+1}\right)=3 n+1$. It can be easily verified that $f$ is an absolutely harmonious labeling of $G_{n}$ and hence $K_{4}$-snakes are absolutely harmonious. For example, an absolutely harmonious labeling of a $K_{4}$-snake with five blocks is shown in Fig.3.7.


Fig.3.7
A quadrilateral snake is obtained from a path $u_{1} u_{2} \ldots u_{n+1}$ by joining $u_{i}, u_{i+1}$ to new vertices $v_{i}, w_{i}$ respectively and joining $v_{i}$ and $w_{i}$.

Theorem 3.8 All quadrilateral snakes are absolutely harmonious.
Proof Let $G_{n}$ be a quadrilateral snake with $V\left(G_{n}\right)=\left\{u_{i}: 1 \leq i \leq n+1\right\} \cup\left\{v_{i}, w_{i}: 1 \leq i \leq n\right\}$ and $E\left(G_{n}\right)=\left\{u_{i} u_{i+1}, u_{i} v_{i}, u_{i+1} w_{i}, v_{i} w_{i}: 1 \leq i \leq n\right\}$. Then $p=3 n+1$ and $q=4 n$. Let $m=\left\{\begin{array}{ll}\frac{n}{2} & \text { if } n \equiv 0(\bmod 2) \\ \frac{n-1}{2} & \text { if } n \equiv 1(\bmod 2)\end{array}\right.$.

Define $f: V\left(G_{n}\right) \longrightarrow\{0,1,2, \ldots 4 n-1\}$ as follows:

$$
\begin{gathered}
f\left(u_{i}\right)=\left\{\begin{array}{ll}
0 & \text { if } i=1 \\
4 i-6 & \text { if } 2 \leq i \leq m+1, \\
4 i-7 & \text { if } m+2 \leq i \leq n+1
\end{array} \quad, \quad f\left(v_{i}\right)= \begin{cases}4 i-3 & \text { if } 1 \leq i \leq m \\
4 i-2 & \text { if } m+1 \leq i \leq n\end{cases} \right. \\
f\left(w_{i}\right)= \begin{cases}4 i & \text { if } 1 \leq i \leq m \\
4 i-1 & \text { if } m+1 \leq i \leq n\end{cases}
\end{gathered}
$$

It can be easily verified that $f$ is an absolutely harmonious labeling of the quadrilateral snake $G_{n}$ and hence quadrilateral snakes are absolutely harmonious. For example, an absolutely harmonious labeling of a quadrilateral snake with six blocks is shown in Fig.3.8.


Fig. 3.8

Theorem 3.9 The disjoint union of $m$ copies of the complete graph on four vertices, $m K_{4}$ is absolutely harmonious.

Proof Let $u_{i}^{j}$ where $1 \leq i \leq 4$ and $1 \leq j \leq m$ denotes the $i^{\text {th }}$ vertex of the $j^{t h}$ copy of $m K_{4}$. We note that that $m K_{4}$ has order $4 m$ and size $6 m$. Define $f: V\left(m K_{4}\right) \longrightarrow\{0,1,2, \ldots 6 m-1\}$ as follows: $f\left(u_{1}^{1}\right)=0, f\left(u_{2}^{1}\right)=1, f\left(u_{3}^{1}\right)=2, f\left(u_{4}^{1}\right)=4, f\left(u_{1}^{2}\right)=q-3, f\left(u_{2}^{2}\right)=q-4, f\left(u_{3}^{2}\right)=$ $q-5, f\left(u_{4}^{2}\right)=q-7, f\left(u_{i}^{j+2}\right)=f\left(u_{i}^{j}\right)+6$ if $j$ is odd, and $f\left(u_{i}^{j+2}\right)=f\left(u_{i}^{j}\right)-6$ if $j$ is even, where $1 \leq i \leq 4$ and $1 \leq j \leq m-2$. Clearly $f$ is an absolutely harmonious labeling. For example, an absolutely harmonious labeling of $5 K_{4}$ is shown in Figure 11.

Box
Observation 3.10 If $f$ is an absolutely harmonious labeling of a graph $G$, which is not a tree, then

1. Each $x$ in the set $\{0,1,2\}$ has inverse image.
2. Inverse images of 0 and 1 are adjacent in $G$.
3. Inverse images of 0 and 2 are adjacent in $G$.

Theorem 3.11 The disjoint union of $m$ copies of the complete graph on three vertices, $m K_{3}$ is absolutely harmonious if and only if $m=1$.

Proof Let $u_{i}^{j}$, where $1 \leq i \leq 3$ and $1 \leq j \leq m$ denote the $i^{\text {th }}$ vertex of the $j^{\text {th }}$ copy of $m K_{3}$. Assignments of the values $0,1,2$ to the vertices of $K_{3}$ gives the desired absolutely harmonious labeling of $K_{3}$. For $m \geq 2, m K_{3}$ has $3 m$ vertices and $3 m$ edges. If $m K_{3}$ is an absolutely harmonious graph, we can assign the numbers $\{0,1,2,3 m-1\}$ to the vertices of $m K_{3}$ in such a way that its edges receive each of the numbers $a_{0}, a_{1}, \ldots, a_{q-1}$ where $a_{i}=$ $q-i$ or $q+i, 0 \leq i \leq q-1$. By Observation 3.10, we can assume, without loss of generality that $f\left(u_{1}^{1}\right)=0, f\left(u_{2}^{1}\right)=1, f\left(u_{3}^{1}\right)=2$. Thus we get the edge labels $a_{q-1}, a_{q-2}$ and $a_{q-3}$. In order to have an edge labeled $a_{q-4}$, we must have two adjacent vertices labeled $q-1$ and $q-3$. we can assume without loss of generality that $f\left(u_{1}^{2}\right)=q-1$ and $f\left(u_{2}^{2}\right)=q-3$. In order to have an edge labeled $a_{q-5}$, we must have $f\left(u_{2}^{3}\right)=q-4$. There is now no way to obtain an edge labeled $a_{q-6}$. This contradiction proves the theorem.

Theorem 3.12 A complete graph $K_{n}$ is absolutely harmonious graph if and only if $n=3$ or 4 .
Proof From the definition of absolutely harmonious labeling, it can be easily verified that $K_{1}$ and $K_{2}$ are not absolutely harmonious graphs. Assignments of the values $0,1,2$ and $0,1,2,4$ respectively to the vertices of $K_{3}$ and $K_{4}$ give the desired absolutely harmonious labeling of them. For $n>4$, the graph $K_{n}$ has $q \geq 10$ edges. If $K_{n}$ is an absolutely harmonious graph, we can assign a subset of the numbers $\{0,1,2, q-1\}$ to the vertices of $K_{n}$ in such a way that the edges receive each of the numbers $a_{0}, a_{1}, \ldots, a_{q-1}$ where $a_{i}=q-i$ or $q+i, 0 \leq i \leq q-1$. By Observation $3.10,0,1$, and 2 must be vertex labels. With vertices labeled 0,1 , and 2 , we have edges labeled $a_{q-1}, a_{q-2}$ and $a_{q-3}$. To have an edge labeled $a_{q-4}$ we must adjoin the vertex label 4. Had we adjoined the vertex label 3 to induce $a_{q-4}$, we would have two edges labeled $a_{q-3}$, namely, between 0 and 3 , and between 1 and 2 . Had we adjoined the vertex labels $q-1$
and $q-3$ to induce $a_{q-4}$, we would have three edges labeled $a_{1}$, namely, between $q-1$ and 0 , between $q-1$ and 2 , and between $q-3$ and 2 . With vertices labeled $0,1,2$, and 4 , we have edges labeled $a_{q-1}, a_{q-2}, a_{q-3}, a_{q-4}, a_{q-5}$, and $a_{q-6}$. Note that for $K_{4}$ with $q=6$, this gives the absolutely harmonious labeling. To have an edge labeled $a_{q-7}$, we must adjoin the vertex label 7 ; all the other choices are ruled out. With vertices labeled $0,1,2,4$ and 7 , we have edges labeled $a_{q-1}, a_{q-2}, a_{q-3}, a_{q-4}, a_{q-5}, a_{q-6}, a_{q-7}, a_{q-8}, a_{q-9}$, and $a_{q-11}$. There is now no way to obtain an edge labeled $a_{q-10}$, because each of the ways to induce $a_{q-10}$ using two numbers contains at least one number that can not be assigned as vertex label. We may easily verify that the following boxed numbers are not possible choices as vertex labels:

| 0 | $\boxed{10}$ |
| ---: | ---: |
| 1 |  |
| 2 | 8 |
| 3 | 7 |
| 4 | 6 |
| $q-1$ | $q-9$ |
| $q-2$ | $q-8$ |
| $q-3$ | $q-7$ |
| $q-4$ | $q-6$ |

This contradiction proves the theorem.

## References

[1] G. J. Chang, D. F. Hsu, and D. G. Rogerss, Additive variations on a graceful theme: Some results on harmonious and other related graphs, Congressus Numerantium, 32 (1981) 181197.
[2] J. A. Gallian, A dynamic survey of graph labeling, The Electronic Journal of Combinatorics, 16 (2009), \#DS6.
[3] R. L. Graham and N. J. A. Solane, On additive bases and Harmonious Graphs, SIAM, J. Alg. Discrete Methods, 1(1980) 382-404.
[4] S. M. Lee, E. Schmeichel, and S. C. Shee, On felicitous graphs, Discrete Mathematics, 93 (1991) 201-209.
[5] Thom Grace, On sequential labelings of graphs, Journal of Graph Theory, 7 (1983), 195 201.

# The Toroidal Crossing Number of $K_{4, n}$ 

Shengxiang Lv<br>(Department of Mathematics, Hunan University of Science and Technology, Xiangtan 411201, China)

Tang Ling, Yuanqiu Huang

(Department of Mathematics, Hunan Normal University, Changsha 410081, China)

E-mail: lsxx23@yahoo.com.cn


#### Abstract

In this paper, we study the crossing number of the complete bipartite graph $K_{4, n}$ in torus and obtain $$
c r_{T}\left(K_{4, n}\right)=\left\lfloor\frac{n}{4}\right\rfloor\left(2 n-4\left(1+\left\lfloor\frac{n}{4}\right\rfloor\right)\right)
$$

Key Words: Smarandache $\mathscr{P}$-drawing, crossing number, complete bipartite graph, torus. AMS(2010): 05C10


## §1. Introduction

A complete bipartite graph $K_{m, n}$ is a graph with vertex set $V_{1} \cup V_{2}$, where $V_{1} \cap V_{2}=\emptyset,\left|V_{1}\right|=m$ and $\left|V_{2}\right|=n$; and with edge set of all pairs of vertices with one element in $V_{1}$ and the other in $V_{2}$. The vertices in $V_{1}$ will be denoted by $b_{i}, b_{j}, b_{k}, \cdots$ and the vertices in $V_{2}$ will be denoted by $a_{i}, a_{j}, a_{k}, \cdots$.

A drawing is a mapping of a graph $G$ into a surface. A Smarandache $\mathscr{P}$-drawing of a graph $G$ for a graphical property $\mathscr{P}$ is such a good drawing of $G$ on the plane with minimal intersections for its each subgraph $H \in \mathscr{P}$. A Smarandache $\mathscr{P}$-drawing is said to be optimal if $\mathscr{P}=G$ and it minimizes the number of crossings. Particularly, a drawing is good if it satisfies: (1) no two arcs which are incident with a common node have a common point; (2) no arc has a self-intersection; (3) no two arcs have more than one point in common; (4) no three arcs have a point in common. A common point of two arcs is called as a crossing. An optimal drawing in a given surface is a good drawing which has the smallest possible number of crossings. This number is the crossing number of the graph in the surface. We denote the crossing number of $G$ in $T$, the torus, by $c r_{T}(G)$, a drawing of $G$ in $T$ by $D$. In this paper, we often speak of the nodes as vertices and the arcs as edges. For more graph terminologies and notations not mentioned here, you can refer to $[1,3]$.

Garey and Johnson [2] stated that determining the crossing number of an arbitrary graph

[^5]is NP-complete. In 1969, Guy and Jenkyns [4] proved that the crossing number of the complete bipartite graph $K_{3, n}$ in torus is $\left\lfloor\frac{(n-3)^{2}}{12}\right\rfloor$, and obtained the bounds on the crossing number of the complete bipartite graph $K_{m, n}$ in torus. In 1971, Kleitman [6] proved that the crossing number of the complete bipartite graph $K_{5, n}$ in plane is $4\left\lfloor\frac{n}{2}\right\rfloor\left\lfloor\frac{n-1}{2}\right\rfloor$ and the crossing number of the complete bipartite graph $K_{6, n}$ in plane is $6\left\lfloor\frac{n}{2}\right\rfloor\left\lfloor\frac{n-1}{2}\right\rfloor$. Later, Richter and S̆irán̆ [7] obtained the crossing number of the complete bipartite graph $K_{3, n}$ in an arbitrary surface. Recently, Ho [5] proved that the crossing number of the complete bipartite graph $K_{4, n}$ in real projective plane is $\left\lfloor\frac{n}{3}\right\rfloor\left(2 n-3\left(1+\left\lfloor\frac{n}{3}\right\rfloor\right)\right)$. In this paper, we obtain the crossing number of the complete bipartite graph $K_{4, n}$ in torus following.

Theorem 1 The crossing number of the complete bipartite graph $K_{4, n}$ in torus is

$$
\operatorname{cr}_{T}\left(K_{4, n}\right)=\left\lfloor\frac{n}{4}\right\rfloor\left(2 n-4\left(1+\left\lfloor\frac{n}{4}\right\rfloor\right)\right)
$$

For convenience, let $f(n)=\left\lfloor\frac{n}{4}\right\rfloor\left(2 n-4\left(1+\left\lfloor\frac{n}{4}\right\rfloor\right)\right)$.

## §2. Some Lemmas

In a drawing $D$ of the complete bipartite $K_{m, n}$ in $T$, we denote by $c r_{D}\left(a_{i}, a_{j}\right)$ the number of crossings on edges one of which is incident with a vertex $a_{i}$ and the other incident with $a_{j}$, and by $c r_{D}\left(a_{i}\right)$ the number of crossings on edges incident with $a_{i}$. Obviously,

$$
c r_{D}\left(a_{i}\right)=\sum_{k=1}^{n} c r_{D}\left(a_{i}, a_{k}\right)
$$

In every good drawing $D$, the crossing number in $D, \operatorname{cr}_{T}(D)$, is

$$
c r_{T}(D)=\sum_{i=1}^{n} \sum_{k=i+1}^{n} c r_{D}\left(a_{i}, a_{k}\right)
$$

$\operatorname{As} c r_{D}\left(a_{i}, a_{i}\right)=0$ for all $i$, hence

$$
\begin{equation*}
c r_{T}(D)=\frac{1}{2} \sum_{i=1}^{n} \sum_{k=1}^{n} c r_{D}\left(a_{i}, a_{k}\right)=\frac{1}{2} \sum_{i=1}^{n} c r_{D}\left(a_{i}\right) . \tag{1}
\end{equation*}
$$



Fig. 1
Fig.1. An optimal drawing of $K_{4,4}$ in $T$

Note that, in a crossing-free drawing of a connected subgraph of the complete bipartite graph $K_{m, n}$, every circuit has an even number of vertices, and in particular, every region into which the edges divide the surface is bounded by an even circuit. So, if $F$ is the number of regions, $E$ the number of edges and $V$ the number of vertices, by the Eular's formula for $T$,

$$
\begin{align*}
& V-E+F \geq 0 \\
& F \geq E-V  \tag{2}\\
& 4 F \leq 2 E \tag{3}
\end{align*}
$$

Suppose we have an optimal drawing of the complete bipartite graph $K_{m, n}$ in $T$, i.e., one with exactly $c r_{T}\left(K_{m, n}\right)$ crossings. Then by deleting $c r_{T}\left(K_{m, n}\right)$ edges, a crossing-free drawing will be obtained. From equations (2) and (3),

$$
E-V=\left(m n-c r_{T}\left(K_{m, n}\right)\right)-(m+n) \leq F \leq \frac{1}{2} E=\frac{1}{2}\left(\left(m n-c r_{T}\left(K_{m, n}\right)\right)\right.
$$

this implies

$$
\begin{equation*}
c r_{T}\left(K_{m, n}\right) \geq m n-2(m+n) \tag{4}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
c r_{T}\left(K_{4, n}\right) \geq 2 n-8 \tag{5}
\end{equation*}
$$

In Fig.1, it is a crossing-free drawing of the complete bipartite graph $K_{4,4}$ in $T$, hence

$$
\begin{equation*}
c r_{T}\left(K_{4,4}\right)=0 \tag{6}
\end{equation*}
$$

In paper [4], the following two lemmas can be find.
Lemma 1 Let $m, n, h$ be positive integers such that the complete bipartite graph $K_{m, h}$ embeds in $T$, then

$$
c r_{T}\left(K_{m, n}\right) \leq \frac{1}{2}\left\lfloor\frac{n}{h}\right\rfloor\left[2 n-h\left(1+\left\lfloor\frac{n}{h}\right\rfloor\right)\right\rfloor\left\lfloor\frac{m}{2}\right\rfloor\left\lfloor\frac{m-1}{2}\right\rfloor .
$$

Lemma 2 If $D$ is a good drawing of the complete bipartite graph $K_{m, n}$ in a surface $\Sigma$ such that, for some $k<n$, some $K_{m, k}$ is optimally drawn in $\Sigma$, then

$$
c r_{\Sigma}(D) \geq c r_{\Sigma}\left(K_{m, k}\right)+(n-k)\left(c r_{\Sigma}\left(K_{m, k+1}\right)-c r_{\Sigma}\left(K_{m, k}\right)\right)+c r_{\Sigma}\left(K_{m, n-k}\right)
$$



Fig. 2


Fig. 3

Lemma 3 For $n \geq 4, c r_{T}\left(K_{4, n}\right) \leq f(n)$; especially, when $4 \leq n \leq 8, c r_{T}\left(K_{4, n}\right)=f(n)$.
Proof As $c r_{T}\left(K_{4,4}\right)=0$, by applying Lemma 1 with $m=h=4$, then $c r_{T}\left(K_{4, n}\right) \leq$ $f(n), n \geq 4$. Especially, as $f(n)=2 n-8$ for $4 \leq n \leq 8$, combining with equation (5), then $c r_{T}\left(K_{4, n}\right)=f(n)$ for $4 \leq n \leq 8$.


Lemma 4 There is no good drawing $D$ of $K_{4,5}$ in $T$ such that
(1) $c r_{D}\left(a_{1}, a_{2}\right)=c r_{D}\left(a_{1}, a_{i}\right)=c r_{D}\left(a_{2}, a_{i}\right)=0$ for $3 \leq i \leq 5$;
(2) $c r_{D}\left(a_{3}, a_{4}\right)=c r_{D}\left(a_{3}, a_{5}\right)=c r_{D}\left(a_{4}, a_{5}\right)=1$.


Proof Note that $T$ can be viewed as a rectangle with its opposite sides identified. As $D$ is a good drawing, by deformation of the edges without changing the crossings and renaming the vertices if necessary, we can assume that the edges incident with $a_{1}$ are drawn as in Fig.2. Since $c r_{D}\left(a_{1}, a_{2}\right)=0$, by deformation of edges without changing the crossings, we also assume that the edge $a_{2} b_{1}$ is drawn as in Fig.3. If the other three edges incident with $a_{2}$ are drawn without passing the sides of the rectangle (see Fig.3), then no matter which region $a_{3}$ is located, we have $\operatorname{cr}_{D}\left(a_{1}, a_{3}\right) \geq 1$ or $\operatorname{cr}_{D}\left(a_{2}, a_{3}\right) \geq 1$.


Fig.7(2)

So, there is at least one edge incident with $a_{2}$ which passes the sides of the rectangle. By deformation without changing the crossings and renaming the vertices if necessary, we assume that edge $a_{2} b_{2}$ passes the top and bottom sides of the rectangle only one time and is drawn as in Fig.4. Then we cut $T$ along the circuit $a_{1} b_{1} a_{2} b_{2} a_{1}$ and obtain a surface which is homeomorphic to a ring in plane, denote by $P$, see Fig.5. Now, we put the vertices $b_{3}, b_{4}$ in $P$ and use two rectangles to represent the outer and inner boundary which are both the circuit $a_{1} b_{1} a_{2} b_{2} a_{1}$.

As the vertices $b_{3}$ and $b_{4}$ are connected to $a_{1}$ and $a_{2}$ either in the outer or in the inner rectangle, which together presents 16 possibilities. In some cases, the four edges can either separate the two rectangles or not, implying up to 32 cases. Using symmetry, several cases are eliminated: without loss of generality, the vertex $b_{3}$ is connected to $a_{2}$ in the outer rectangle.


Fig.8(1)


Fig.8(3)


Fig.8(2)


Fig.8(4)

First, assume that $b_{3}$ is also connected to $a_{1}$ in the outer rectangle. If $b_{4}$ is connected to both $a_{1}$ and $a_{2}$ in the outer rectangle, we obtain Fig.6(1) if the four edges separate the two rectangles, and Fig.6(2) if they do not. If $b_{4}$ is connected to $a_{1}$ in the inner rectangle and $a_{2}$ in the outer rectangle, we obtain Fig.6(3). If it is connected to $a_{1}$ in the outer rectangle and $a_{2}$ in the inner rectangle, then by relabeling $a_{1}$ and $a_{2}$, we obtain Fig.6(3). If $b_{4}$ is connected to both $a_{1}$ and $a_{2}$ in the inner rectangle, we obtain Fig.6(4).

Second, assume that $b_{3}$ is connected to $a_{1}$ in the inner rectangle. If $b_{4}$ is connected to both $a_{1}$ and $a_{2}$ in the outer rectangle, then by relabeling of $b_{3}$ and $b_{4}$, we obtain Fig.6(3). If $b_{4}$ is connected to $a_{1}$ in the inner rectangle and $a_{2}$ in the outer rectangle, we obtain Fig.6(5) if the four edges separate the two rectangles, and Fig.6(6) if they do not. If $b_{4}$ is connected to $a_{2}$ in the inner rectangle and $a_{1}$ in the outer rectangle, we obtain Fig.6(7). Finally, if $b_{4}$ is connected to both $a_{1}$ and $a_{2}$ in the inner rectangle, we obtain Fig.6(8).


Now, by drawing Fig.6(1) back into $T$ and cut $T$ along the circuit $a_{1} b_{2} a_{2} b_{4} a_{1}$, we obtain Fig.7(1); by drawing Fig.6(6) back into $T$ and cut $T$ along the circuit $a_{1} b_{4} a_{2} b_{2} a_{1}$, we obtain Fig.7(2). It is easy to find out that Fig.7(1) and Fig.6(4), Fig.7(2) and Fig.6(3) have the same structure if ignoring the labels of $b$. In Fig.6(8), by exchanging the inner and outer rectangles and the labels of $b_{3}, b_{4}$, we obtain Fig.6(3). In Fig.6(2), as each region has at most 3 vertices of $\left\{b_{1}, b_{2}, b_{3}, b_{4}\right\}$ on its boundary, we will have $c r_{D}\left(a_{1}, a_{i}\right) \geq 1$ or $c r_{D}\left(a_{2}, a_{i}\right) \geq 1$ for $i=3,4,5$. So, we only need to consider the cases in Fig.6(3-5,7).

In Fig.6(3), since $c r_{D}\left(a_{1}, a_{3}\right)=c r_{D}\left(a_{2}, a_{3}\right)=0$, we can draw the edges incident with $a_{3}$ in four different ways, see Fig.8(1-4). Furthermore, as $c r_{D}\left(a_{1}, a_{4}\right)=c r_{D}\left(a_{2}, a_{4}\right)=0$ and $\operatorname{cr}_{D}\left(a_{3}, a_{4}\right)=1, a_{4}$ can only be putted in region I or II. In Fig.8(3-4), we can draw the edges incident with $a_{4}$ in four different ways, see Fig.9(1-4). In Fig.8(1-2), there are also four different ways to draw the edges incident with $a_{4}$, but they can be obtained by relabeling $a_{3}$ and $a_{4}$ in Fig.9((1-4). Then, we can see that no matter which region $a_{5}$ lies, we cannot have $c r_{D}\left(a_{3}, a_{5}\right)=$ $c r_{D}\left(a_{4}, a_{5}\right)=1$.


Fig.10(1)


Fig.10(2)


Fig.10(3)

In Fig.6(4), we have only one way to draw the edges incident with $a_{3}$, see Fig.10(1). Furthermore, we have two drawings of $a_{4}$ in Fig.10(1), see Fig.10(2-3). But, by observation, we
cannot have $c r_{D}\left(a_{3}, a_{5}\right)=c r_{D}\left(a_{4}, a_{5}\right)=1$.
In Fig.6(5,7), no matter which regions $a_{3}, a_{4}$ locate, we will have $\operatorname{cr}_{D}\left(a_{3}, a_{4}\right) \geq 2$ or $c r_{D}\left(a_{3}, a_{4}\right)=0$. Now, the proof completes.

## §3. The proof of the Main Theorem

The proof of Theorem 1 is by induction on $n$. The base of the induction is $n \leq 8$ and has been obtained from Lemma 3. For $n \geq 9$, by Lemma 3, we only need to prove that $c r_{T}\left(K_{4, n}\right) \geq f(n)$. Let $n=4 q+r$ where $0 \leq r \leq 3$, and $D$ be an optimal drawing of $K_{4, n}$ in $T$.

First, we assume that there exists a $K_{4,4}$ in $D$ which is drawn without crossings. From Lemma $3, c r_{T}\left(K_{4,5}\right)=2$, and by the inductive assumption, $c r_{T}\left(K_{4, n-4}\right)=f(n-4)$. Hence, by applying Lemma 2 with $m=k=4$,

$$
\begin{aligned}
c r_{T}(D) & \geq 2(n-4)+f(n-4)=2(n-4)+\left\lfloor\frac{n-4}{4}\right\rfloor\left(2(n-4)-4\left(1+\left\lfloor\frac{n-4}{4}\right\rfloor\right)\right) \\
& =8 q+2 r-8+(q-1)(4 q+2 r-8)=4 q^{2}+2 q r-4 q
\end{aligned}
$$

which is $f(n)$, since

$$
\begin{equation*}
f(n)=\left\lfloor\frac{n}{4}\right\rfloor\left(2 n-4\left(1+\left\lfloor\frac{n}{4}\right\rfloor\right)\right)=q(8 q+2 r-4(1+q))=4 q^{2}+2 q r-4 q \tag{7}
\end{equation*}
$$

Second, we assume that every $K_{4,4}$ in $D$ is drawn with at least one crossings. Clearly, $K_{4, n}$ contains $n$ subgraphs $K_{4, n-1}$, each contains at least $f(n-1)$ crossings by the inductive hypothesis. As each crossing will be counted $n-2$ times, hence

$$
\begin{equation*}
c r_{T}(D) \geq \frac{n}{n-2} c r_{T}\left(K_{4, n-1}\right)=\frac{n}{n-2} f(n-1) \tag{8}
\end{equation*}
$$

From equation (7),

$$
f(n)= \begin{cases}q(4 q-4), & \text { for } n=4 q \\ q(4 q-2), & \text { for } n=4 q+1 \\ 4 q^{2}, & \text { for } n=4 q+2 \\ q(4 q+2), & \text { for } n=4 q+3\end{cases}
$$

Combining this with equation (8),

$$
c_{T}(D) \geq \begin{cases}q(4 q-4), & \text { for } n=4 q \\ q(4 q-2)-1-\frac{2 q+1}{4 q-1}, & \text { for } n=4 q+1 \\ 4 q^{2}-1, & \text { for } n=4 q+2 \\ q(4 q+2)-\frac{2 q}{4 q+1}, & \text { for } n=4 q+3\end{cases}
$$

As $n \geq 9$, namely $q \geq 2$, and the crossing number is an integer, thus, when $n=4 q$ or $4 q+3$,

$$
c r_{T}\left(K_{4, n}\right)=c r_{T}(D) \geq f(n)
$$

when $n=4 q+1$ or $4 q+2$,

$$
\operatorname{cr}_{T}\left(K_{4, n}\right)=\operatorname{cr}_{T}(D) \geq f(n)-1
$$

Therefore, only the two cases $n=4 q+1$ and $n=4 q+2$ are needed considering. In the following, we assume that $c r_{T}\left(K_{4, n}\right)=c r_{T}(D)=f(n)-1$ for $n=4 q+1$ or $4 q+2$, and denote the drawing of $K_{4, n-1}$ obtained by deleting the vertex $a_{i}$ of $K_{4, n}$ in $D$ by $D-\left\{a_{i}\right\}$.

Case 1. $n=4 q+1$.
By the inductive assumption,

$$
\operatorname{cr}_{T}\left(D-\left\{a_{i}\right\}\right) \geq f(4 q), 1 \leq i \leq 4 q+1
$$

As $c r_{T}(D)=f(4 q+1)-1=4 q^{2}-2 q-1$, then

$$
c r_{D}\left(a_{i}\right)=c r_{T}(D)-c r_{T}\left(D-\left\{a_{i}\right\}\right) \leq f(4 q+1)-1-f(4 q)=2 q-1,1 \leq i \leq 4 q+1
$$

Let $x$ be the number of $a_{i}$ such that $\operatorname{cr}_{D}\left(a_{i}\right)=2 q-1, y$ be the number of $a_{i}$ such that $c r_{D}\left(a_{i}\right)=2 q-2$, thus, the number of $a_{i}$ such that $c r_{D}\left(a_{i}\right) \leq 2 q-3$ is $4 q+1-(x+y)$. By equation(1), it holds

$$
\begin{aligned}
(2 q-1) x & +(2 q-2) y+(4 q+1-x-y)(2 q-3) \geq 2 c r_{T}(D)=8 q^{2}-4 q-2 \\
2 x+y & \geq 6 q+1
\end{aligned}
$$

As $x+y \leq 4 q+1$, then $x \geq 2 q$. Without loss of generality, by renaming the vertices, suppose that $c r_{D}\left(a_{i}\right)=2 q-1$ for $i \leq x$.

Case 1.1 There exists a pair of $(i, j), 1 \leq i<j \leq x$, such that $c r_{D}\left(a_{i}, a_{j}\right)=0$. Denote the drawing of the graph $K_{4,4 q-1}$ obtained by deleting the vertices $a_{i}, a_{j}$ of the graph $K_{4,4 q+1}$ in $D$ by $D-\left\{a_{i}, a_{j}\right\}$. Then,

$$
\operatorname{cr}_{T}\left(D-\left\{a_{i}, a_{j}\right\}\right)=f(4 q+1)-1-2(2 q-1)=4 q^{2}-6 q+1
$$

But this contradicts the inductive assumption that $c r_{T}\left(K_{4,4 q-1}\right)=f(4 q-1)=4 q^{2}-6 q+2$.
Case 1.2 For every $(i, j), 1 \leq i<j \leq x, \operatorname{cr}_{D}\left(a_{i}, a_{j}\right) \geq 1$. As $c r_{D}\left(a_{i}\right)=2 q-1$, obviously, $x=2 q$ and

$$
c r_{D}\left(a_{i}, a_{j}\right)=1,1 \leq i<j \leq 2 q, c r_{D}\left(a_{i}, a_{h}\right)=0,1 \leq i \leq 2 q<h \leq 4 q+1
$$

Furthermore, as $x+y \leq 4 q+1$ and $2 x+y \geq 6 q+1$, then $y=2 q+1$. By the definition of $y$, there exist $a_{h}, a_{k}$, where $2 q+1 \leq h<k \leq 4 q+1$, such that $c r_{D}\left(a_{h}, a_{k}\right)=0$. Now, we obtain a drawing of $K_{4,5}$ in $T$ with vertices $a_{h}, a_{k}, a_{1}, a_{2}, a_{3}$ such that $\operatorname{cr}_{D}\left(a_{h}, a_{k}\right)=\operatorname{cr}_{D}\left(a_{h}, a_{i}\right)=$ $c r_{D}\left(a_{k}, a_{i}\right)=0(1 \leq i \leq 3)$ and $c r_{D}\left(a_{1}, a_{2}\right)=c r_{D}\left(a_{1}, a_{3}\right)=c r_{D}\left(a_{2}, a_{3}\right)=1$. Contradicts to Lemma 4.

Combining the above two subcases, we have $c r_{T}\left(K_{4,4 q+1}\right)=f(4 q+1)=q(4 q-2)$.
Case 2. $n=4 q+2$.
By the inductive assumption,

$$
\operatorname{cr}_{T}\left(D-\left\{a_{i}\right\}\right) \geq f(4 q+1)=q(4 q-2), 1 \leq i \leq 4 q+2
$$

As $c r_{T}(D)=f(4 q+2)-1=4 q^{2}-1$, thus

$$
c r_{D}\left(a_{i}\right)=c r_{T}(D)-c r_{T}\left(D-\left\{a_{i}\right\}\right) \leq(f(4 q+2)-1)-f(4 q+1)=2 q-1
$$

Let $t$ be the number of $a_{i}$ such that $\operatorname{cr}_{D}\left(a_{i}\right)=2 q-1$, then there are $(4 q+2-t)$ vertices $a_{i}$ such that $c r_{D}\left(a_{i}\right) \leq 2 q-2$. From equation (1),

$$
\begin{aligned}
(2 q-1) t & +(2 q-2)(4 q+2-t) \geq 2 c r_{T}(D)=8 q^{2}-2 \\
t & \geq 4 q+2
\end{aligned}
$$

As $t \leq n=4 q+2$, hence, $t=4 q+2$, this implies that $c r_{D}\left(a_{i}\right)=2 q-1(1 \leq i \leq 4 q+2)$.
If there exists a pair of $(i, j), 1 \leq i<j \leq 4 q+2$, such that $c r_{D}\left(a_{i}, a_{j}\right) \geq 3$, then,

$$
c r_{T}\left(D-\left\{a_{i}\right\}\right)=c r_{T}(D)-c r_{D}\left(a_{i}\right)=4 q^{2}-1-(2 q-1)=4 q^{2}-2 q
$$

and

$$
\operatorname{cr}_{\left(D-\left\{a_{i}\right\}\right)}\left(a_{j}\right)=c r_{D}\left(a_{j}\right)-c r_{D}\left(a_{i}, a_{j}\right) \leq 2 q-1-3=2 q-4
$$

Now, by putting a new vertex $a_{i}^{\prime}$ near the vertex $a_{j}$ in $D-\left\{a_{i}\right\}$ and drawing the edges $a_{i}^{\prime} b_{k}(1 \leq$ $k \leq 4)$ nearly to $a_{j} b_{k}$, a new drawing of $K_{4,4 q+2}$ in $T$ is obtained, denoted by $D^{\prime}$. Clearly,

$$
c r_{D^{\prime}}\left(a_{i}^{\prime}, a_{j}\right)=2 \text { and } c r_{D^{\prime}}\left(a_{i}^{\prime}, a_{h}\right)=c r_{D-\left\{a_{i}\right\}}\left(a_{j}, a_{h}\right), h \neq j .
$$

Thus,

$$
\operatorname{cr}_{T}\left(D^{\prime}\right)=c r_{T}\left(D-\left\{a_{i}\right\}\right)+2+c r_{\left(D-\left\{a_{i}\right\}\right)}\left(a_{j}\right) \leq 4 q^{2}-2
$$

But, this contradicts to the hypothesis that $c r_{T}\left(K_{4,4 q+2}\right) \geq 4 q^{2}-1$.
Therefore, for $1 \leq i<j \leq 4 q+2, c r_{D}\left(a_{i}, a_{j}\right) \leq 2$. For each $a_{i}, 1 \leq i \leq 4 q+2$, let

$$
\begin{aligned}
S_{0}^{(i)}=\left\{a_{j} \mid c r_{D}\left(a_{i}, a_{j}\right)=0, j \neq i\right\}, & S_{\geq 1}^{(i)}=\left\{a_{j} \mid c r_{D}\left(a_{i}, a_{j}\right) \geq 1\right\} \\
S_{1}^{(i)}=\left\{a_{j} \mid c r_{D}\left(a_{i}, a_{j}\right)=1\right\}, & S_{2}^{(i)}=\left\{a_{j} \mid c r_{D}\left(a_{i}, a_{j}\right)=2\right\}
\end{aligned}
$$

As $c r_{D}\left(a_{i}, a_{j}\right) \leq 2, c r_{D}\left(a_{i}\right)=2 q-1$ is odd, then, for $1 \leq i \leq 4 q+2$,

$$
\begin{equation*}
\emptyset \neq S_{1}^{(i)} \subseteq S_{\geq 1}^{(i)}, \quad\left|S_{1}^{(i)}\right|+\left|S_{2}^{(i)}\right|=\left|S_{\geq 1}^{(i)}\right|, \quad\left|S_{\geq 1}^{(i)}\right|=2 q-1-\left|S_{2}^{(i)}\right| \tag{9}
\end{equation*}
$$

Furthermore, since $q \geq 2$,

$$
\left|S_{0}^{(i)}\right|=4 q+2-1-\left|S_{\geq 1}^{(i)}\right|=2 q+2+\left|S_{2}^{(i)}\right| \geq 6
$$

For $1 \leq i<j \leq 4 q+2$, clearly,

$$
S_{0}^{(i)} \cup S_{\geq 1}^{(i)} \cup\left\{a_{i}\right\}=S_{0}^{(j)} \cup S_{\geq 1}^{(j)} \cup\left\{a_{j}\right\}
$$

If $c r_{D}\left(a_{i}, a_{j}\right)=0$ and $S_{\geq 1}^{(i)} \cap S_{\geq 1}^{(j)}=\emptyset$, then, the above equation implies that

$$
\begin{equation*}
S_{\geq 1}^{(i)} \subseteq S_{0}^{(j)} \quad \text { and } \quad S_{\geq 1}^{(j)} \subseteq S_{0}^{(i)} \tag{10}
\end{equation*}
$$

Without loss of generality, let

$$
\left|S_{2}^{(1)}\right|=\max \left\{\left|S_{2}^{(i)}\right| \mid 1 \leq i \leq 4 q+2\right\}, \quad\left|S_{2}^{(2)}\right|=\max \left\{\left|S_{2}^{(j)}\right| \mid a_{j} \in S_{0}^{(1)}\right\}
$$

For $3 \leq i \leq 4 q+2$, if $a_{i} \notin S_{\geq 1}^{(1)} \cup S_{\geq 1}^{(2)}$, then $a_{i} \in S_{0}^{(1)} \cap S_{0}^{(2)}$. This means that

$$
\left|S_{0}^{(1)} \cap S_{0}^{(2)}\right|=4 q-\left|S_{\geq 1}^{(1)} \cup S_{\geq 1}^{(2)}\right|=4 q-\left|S_{\geq 1}^{(1)}\right|-\left|S_{\geq 1}^{(2)}\right|+\left|S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)}\right|
$$

From equation (9), then

$$
\begin{equation*}
\left|S_{0}^{(1)} \cap S_{0}^{(2)}\right|=2+\left|S_{2}^{(1)}\right|+\left|S_{2}^{(2)}\right|+\left|S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)}\right| \tag{11}
\end{equation*}
$$

With these notations, it is obvious that $\left|S_{2}^{(1)}\right| \geq\left|S_{2}^{(2)}\right|$ and $\operatorname{cr}_{D}\left(a_{1}, a_{2}\right)=0$. In the following, the discussions are divided into two subcases according to $S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)}=\emptyset$ or not.
Case 2.1 $\quad S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)} \neq \emptyset$. Let $\left|S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)}\right|=\alpha \geq 1$, from equation (11),

$$
\left|S_{0}^{(1)} \cap S_{0}^{(2)}\right|=2+\left|S_{2}^{(1)}\right|+\left|S_{2}^{(2)}\right|+\alpha
$$

First, we choose a vertex from $S_{0}^{(1)} \cap S_{0}^{(2)}$, without loss of generality, denoted by $a_{3}$. By the assumption that every $K_{4,4}$ in $D$ is drawn with at least one crossings, hence $c r_{D}\left(a_{3}, a_{i}\right) \geq 1$ for all $a_{i} \in S_{0}^{(1)} \cap S_{0}^{(2)}, a_{i} \neq a_{3}$. Let $U=\left\{a_{i} \mid c r_{D}\left(a_{3}, a_{i}\right)=1, a_{i} \in S_{0}^{(1)} \cap S_{0}^{(2)}\right\}$. Since $a_{3} \in S_{0}^{(1)}$ and $\left|S_{2}^{(2)}\right|=\max \left\{\left|S_{2}^{(j)}\right| \mid a_{j} \in S_{0}^{(1)}\right\}$, then $\left|S_{2}^{(3)}\right| \leq\left|S_{2}^{(2)}\right|$ and

$$
|U| \geq\left|S_{0}^{(1)} \cap S_{0}^{(2)}\right|-1-\left|S_{2}^{(3)}\right| \geq 1+\left|S_{2}^{(1)}\right|+\alpha
$$

Second, we choose a vertex from $U$, denoted by $a_{4}$. By the assumption that every $K_{4,4}$ in $D$ is drawn with at least one crossings, $\operatorname{cr}_{D}\left(a_{4}, a_{i}\right) \geq 1$ for all $a_{i} \in U, a_{i} \neq a_{4}$. As $\left|S_{2}^{(4)}\right| \leq\left|S_{2}^{(1)}\right|$ (for $\left.\left|S_{2}^{(1)}\right|=\max \left\{\left|S_{2}^{(i)}\right| \mid 1 \leq i \leq 4 q+2\right\}\right)$, thus $\left|U \backslash S_{2}^{(4)}\right| \geq \alpha \geq 1$ and there exists one vertex in $U$, denoted by $a_{5}$, such that $\operatorname{cr}_{D}\left(a_{4}, a_{5}\right)=1$. Now, we have a drawing of $K_{4,5}$ in $T$ with vertices $a_{1}, a_{2}, a_{3}, a_{4}, a_{5}$ such that $\operatorname{cr}_{D}\left(a_{1}, a_{2}\right)=c r_{D}\left(a_{1}, a_{k}\right)=c r_{D}\left(a_{2}, a_{k}\right)=0$ for $3 \leq k \leq 5$ and $\operatorname{cr}_{D}\left(a_{3}, a_{4}\right)=c r_{D}\left(a_{3}, a_{5}\right)=c r_{D}\left(a_{4}, a_{5}\right)=1$. But, this contradicts to Lemma 4.
Case $2.2 \quad S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)}=\emptyset$. From equation (11),

$$
\left|S_{0}^{(1)} \cap S_{0}^{(2)}\right|=2+\left|S_{2}^{(1)}\right|+\left|S_{2}^{(2)}\right|
$$

We choose a vertex from $S_{0}^{(1)} \cap S_{0}^{(2)}$, also denoted by $a_{3}$. By the same discussion as in case 2.1, we have $c r_{D}\left(a_{3}, a_{i}\right) \geq 1$ for all $a_{i} \in S_{0}^{(1)} \cap S_{0}^{(2)}, a_{i} \neq a_{3}$. Let $\Lambda=\left\{a_{i} \mid c r_{D}\left(a_{3}, a_{i}\right)=2, a_{i} \in\right.$ $\left.S_{0}^{(1)} \cap S_{0}^{(2)}\right\}, \Phi=\left\{a_{i} \mid c r_{D}\left(a_{3}, a_{i}\right)=1, a_{i} \in S_{0}^{(1)} \cap S_{0}^{(2)}\right\}$. As $a_{3} \in S_{0}^{(1)},\left|S_{2}^{(2)}\right|=\max \left\{\left|S_{2}^{(j)}\right| \mid a_{j} \in\right.$ $\left.S_{0}^{(1)}\right\}$ and $\left|S_{2}^{(1)}\right|=\max \left\{\left|S_{2}^{(i)}\right| \mid 1 \leq i \leq 4 q+2\right\}$, then

$$
\begin{equation*}
\Lambda \subseteq S_{2}^{(3)}, \quad|\Lambda| \leq\left|S_{2}^{(3)}\right| \leq\left|S_{2}^{(2)}\right| \leq\left|S_{2}^{(1)}\right| \tag{12}
\end{equation*}
$$

and

$$
\begin{equation*}
|\Phi|=\left|S_{0}^{(1)} \cap S_{0}^{(2)}\right|-1-|\Lambda|=1+\left|S_{2}^{(1)}\right|+\left|S_{2}^{(2)}\right|-|\Lambda| \tag{13}
\end{equation*}
$$

If there are two vertices in $\Phi$, denoted by $a_{4}, a_{5}$, such that $c r_{D}\left(a_{4}, a_{5}\right)=1$. Then we also have a drawing of $K_{4,5}$ with vertices $a_{1}, a_{2}, a_{3}, a_{4}, a_{5}$ which will contradict to Lemma 4. Hence,
for all $a_{i}, a_{j} \in \Phi\left(a_{i} \neq a_{j}\right), c r_{D}\left(a_{i}, a_{j}\right) \neq 1$, this implies that $c r_{D}\left(a_{i}, a_{j}\right)=2$ since $c r_{D}\left(a_{i}, a_{j}\right)$ cannot be zero (otherwise there exists $K_{4,4}$ in $D$ drawn with no crossings), and

$$
\left|S_{2}^{(i)}\right| \geq|\Phi|-1
$$

Furthermore, if $|\Lambda|<\left|S_{2}^{(2)}\right|$, by equation(13), $|\Phi|>1+\left|S_{2}^{(1)}\right|$, and for each $a_{i} \in \Phi$,

$$
\left|S_{2}^{(i)}\right| \geq|\Phi|-1>\left|S_{2}^{(1)}\right|
$$

This contradicts the maximum of $\left|S_{2}^{(1)}\right|$. Thus,

$$
|\Lambda|=\left|S_{2}^{(2)}\right|, \quad|\Phi|=1+\left|S_{2}^{(1)}\right|
$$

and for each $a_{i} \in \Phi$,

$$
\left|S_{2}^{(i)}\right| \geq|\Phi|-1=\left|S_{2}^{(1)}\right|
$$

As $\left|S_{2}^{(i)}\right| \leq\left|S_{2}^{(2)}\right| \leq\left|S_{2}^{(1)}\right|$, combining equation (12),

$$
\begin{equation*}
\left|S_{2}^{(1)}\right|=\left|S_{2}^{(2)}\right|=\left|S_{2}^{(3)}\right|=\left|S_{2}^{(i)}\right| \tag{14}
\end{equation*}
$$

and

$$
S_{2}^{(3)}=\Lambda \subseteq S_{0}^{(1)} \cap S_{0}^{(2)}
$$

Combining equations (14) and (9), for each $a_{i} \in \Phi$,

$$
\left|S_{\geq 1}^{(1)}\right|=\left|S_{\geq 1}^{(2)}\right|=\left|S_{\geq 1}^{(3)}\right|=\left|S_{\geq 1}^{(i)}\right|
$$

and

$$
\left|S_{1}^{(1)}\right|=\left|S_{1}^{(2)}\right|=\left|S_{1}^{(3)}\right|=\left|S_{1}^{(i)}\right|
$$

As $|\Phi|=1+\left|S_{2}^{(1)}\right|+\left|S_{2}^{(2)}\right|-|\Lambda| \geq 1$, we choose a vertex from $\Phi$ and denote it by $a_{4}$.
If there exists a pair of $(i, j), i \in\{1,2\}$ and $j \in\{3,4\}$, such that $S_{>1}^{(i)} \cap S_{>1}^{(j)} \neq \emptyset$, by replacing $S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)} \neq \emptyset$ with $S_{\geq 1}^{(i)} \cap S_{\geq 1}^{(j)} \neq \emptyset$ in case 2.1, as $a_{j} \in S_{0}^{(1)} \cap \bar{S}_{0}^{(2)}(j=3,4)$ and $\left|S_{2}^{(i)}\right|=\left|S_{2}^{(j)}\right|=\max \left\{\left|S_{2}^{(k)}\right| \mid 1 \leq k \leq 4 q+2\right\}$, we also can obtain a contradiction to Lemma 4.

So, for every $(i, j), i \in\{1,2\}$ and $j \in\{3,4\}, S_{\geq 1}^{(i)} \cap S_{\geq 1}^{(j)}=\emptyset$. As $S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)}=\emptyset$ and $c r_{D}\left(a_{i}, a_{j}\right)=c r_{D}\left(a_{1}, a_{2}\right)=0$, combining equations (9) and (10), then

$$
\emptyset \neq S_{1}^{(1)} \subseteq S_{\geq 1}^{(1)} \subseteq S_{0}^{(2)} \cap S_{0}^{(3)} \cap S_{0}^{(4)} \quad \text { and } \quad \emptyset \neq S_{1}^{(2)} \subseteq S_{\geq 1}^{(2)} \subseteq S_{0}^{(1)} \cap S_{0}^{(3)} \cap S_{0}^{(4)}
$$

Since $S_{1}^{(1)} \neq \emptyset$, there exists a vertex, denoted by $a_{5}$, such that $a_{5} \in S_{1}^{(1)} \subseteq S_{0}^{(2)} \cap S_{0}^{(3)} \cap S_{0}^{(4)}$. This implies that

$$
c r_{D}\left(a_{1}, a_{5}\right)=1 \text { and } c r_{D}\left(a_{2}, a_{5}\right)=c r_{D}\left(a_{3}, a_{5}\right)=c r_{D}\left(a_{4}, a_{5}\right)=0 .
$$

As $S_{\geq 1}^{(2)} \cap S_{\geq 1}^{(3)}=\emptyset,\left|S_{2}^{(1)}\right|=\left|S_{2}^{(2)}\right|=\left|S_{2}^{(3)}\right|, c r_{D}\left(a_{2}, a_{3}\right)=0$ and $a_{5} \subseteq S_{0}^{(2)} \cap S_{0}^{(3)}$, by replacing $S_{\geq 1}^{(1)} \cap S_{\geq 1}^{(2)}=\emptyset$ with $S_{\geq 1}^{(2)} \cap S_{>1}^{(3)}=\emptyset$ and replacing $a_{3}$ with $a_{5}$ in the beginning part of Case 2.2, we also can obtain that $\left|S_{1}^{(5)}\right|=\left|S_{1}^{(2)}\right|=\left|S_{1}^{(3)}\right|$ and $S_{2}^{(5)} \subseteq S_{0}^{(2)} \cap S_{0}^{(3)}$. This means that, for any vertex $a_{k} \in S_{1}^{(2)}$,

$$
\begin{equation*}
c r_{D}\left(a_{5}, a_{k}\right) \leq 1 \tag{15}
\end{equation*}
$$

As $S_{1}^{(2)} \neq \emptyset$, there exists one vertex in $S_{1}^{(2)}$, denoted by $a_{6}$, such that $c r_{D}\left(a_{5}, a_{6}\right)=0$. Otherwise, from equation (15) and $\operatorname{cr}_{D}\left(a_{1}, a_{5}\right)=1, S_{1}^{(2)} \cup\left\{a_{1}\right\} \subseteq S_{1}^{(5)}$. As $a_{1} \notin S_{1}^{(2)}$, then $\left|S_{1}^{(5)}\right| \geq\left|S_{1}^{(2)}\right|+1$, which contradicts to $\left|S_{1}^{(5)}\right|=\left|S_{1}^{(2)}\right|=\left|S_{1}^{(3)}\right|$. Furthermore, as $a_{6} \in S_{1}^{(2)} \subseteq$ $S_{0}^{(1)} \cap S_{0}^{(3)} \cap S_{0}^{(4)}$, we also have

$$
\operatorname{cr}_{D}\left(a_{2}, a_{6}\right)=1 \text { and } \operatorname{cr}_{D}\left(a_{1}, a_{6}\right)=c r_{D}\left(a_{3}, a_{6}\right)=\operatorname{cr}_{D}\left(a_{4}, a_{6}\right)=0 .
$$

Hence, we obtain a good drawing of $K_{4,6}$ in $T$, denoted by $D^{\prime}$, with

$$
\operatorname{cr}_{D^{\prime}}\left(a_{i}\right)=\sum_{j=1}^{6} \operatorname{cr}_{D}\left(a_{i}, a_{j}\right)=1,1 \leq i \leq 6,
$$

and

$$
c r_{T}\left(K_{4,6}\right) \leq c r_{T}\left(D^{\prime}\right)=\frac{1}{2} \sum_{i=1}^{6} c r_{D^{\prime}}\left(a_{i}\right)=3 .
$$

This contradicts to Lemma 3. Thus, $c r_{T}\left(K_{4,4 q+2}\right)=c r_{T}(D)=f(4 q+2)=4 q^{2}$.

## References

[1] Bondy J. A., Murty U. S. R., Graph theory with applications, North Holland, New York, 1982.
[2] M. R. Garey and D. S. Johnson, Crossing number is NP-complete, SIAM J. Alg. Disc. Meth., 4(1983) 312-316.
[3] J. L. Gross, Topological Graph Theory, New York: Wiley, 1989.
[4] R. K. Guy, T. A. Jenkyns, The toroidal crossing number of $K_{m, n}$, J. Combin. Theory Ser. B., 6(1969), 235-250.
[5] P. T. Ho, The crossing number of $K_{4, n}$ on the real projective plane, Discrete Math., 304(2005), 23-33.
[6] D. J. Kleitman, The crossing number of $K_{5, n}$, J. Combin. Theory Ser. B., 9(1971), 315323.
[7] R. B. Richter, J. S̆irán̆, The crossing number of $K_{3, n}$ in a surface, J. Graph Theory, 21(1996), 51-54.

# On Pathos Semitotal and Total Block Graph of a Tree 

Muddebihal M. H.<br>(Department of Mathematics, Gulbarga University, Gulbarga, India) Syed Babajan<br>(Department of Mathematics, Ghousia College of Engineering, Ramanagaram, India)<br>E-mail: babajan.ghousia@gmail.com


#### Abstract

In this communications, the concept of pathos semitotal and total block graph of a graph is introduced. Its study is concentrated only on trees. We present a characterization of those graphs whose pathos semitotal block graphs are planar, maximal outer planar, non-minimally non-outer planar, non-Eulerian and hamiltonian. Also, we present a characterization of graphs whose pathos total block graphs are planar, maximal outer planar, minimally non-outer planar, non-Eulerian, hamiltonian and graphs with crossing number one.


Key Words: Pathos, path number, Smarandachely block graph, semitotal block graph, Total block graph, pathos semitotal graph, pathos total block graph, pathos length, pathos point, inner point number.

AMS(2010): 05C75

## §1. Introduction

The concept of pathos of a graph $G$ was introduced by Harary [2], as a collection of minimum number of line disjoint open paths whose union is $G$. The path number of a graph $G$ is the number of paths in pathos. A new concept of a graph valued functions called the semitotal and total block graph of a graph was introduced by Kulli [6]. For a graph $G(p, q)$ if $B=$ $\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{r} ; r \geq 2\right\}$ is a block of $G$, then we say that point $u_{1}$ and block $B$ are incident with each other, as are $u_{2}$ and $B$ and so on. If two distinct blocks $B_{1}$ and $B_{2}$ are incident with a common cut point, then they are adjacent blocks. The points and blocks of a graph are called its members. A Smarandachely block graph $T_{S}^{V}(G)$ for a subset $V \subset V(G)$ is such a graph with vertices $V \cup \mathcal{B}$ in which two points are adjacent if and only if the corresponding members of $G$ are adjacent in $\langle V\rangle_{G}$ or incident in $G$, where $\mathcal{B}$ is the set of blocks of $G$. The semitotal block graph of a graph $G$ denoted by $T_{b}(G)$ is defined as the graph whose point set is the union of set of points, set of blocks of $G$ in which two points are adjacent if and only if members of $G$ are incident, thus a Smarandachely block graph with $V=\emptyset$. The total block graph of a graph

[^6]$G$ denoted by $T_{B}(G)$ is defined as the graph whose point set is the union of set of points, set of blocks of $G$ in which two points are adjacent if and only if the corresponding members of $G$ are adjacent or incident, i.e., a Smarandachely block graph with $V=V(G)$. Stanton [11] and Harary [3] have calculated the path number for certain classes of graphs like trees and complete graphs.

All undefined terminology will conform with that in Harary [1]. All graphs considered here are finite, undirected and without loops or multiple lines.

The pathos semitotal block graph of a tree $T$ denoted by $P_{T_{B}}(T)$ is defined as the graph whose point set is the union of set of points, set of blocks and the set of path of pathos of $T$ in which two points are adjacent if and only if the corresponding members of $G$ are incident and the lines lie on the corresponding path $P_{i}$ of pathos. Since the system of pathos for a tree is not unique, the corresponding pathos semitotal and pathos total block graph of a tree $T$ is also not unique.

In Fig.1, a tree $T$, its semitotal block graph $T_{b}(T)$ and their pathos semi total block $P_{T_{b}}(T)$ graph are shown. In Fig. 2, a tree $T$, its semitotal block graph $T_{b}(T)$ and their pathos total block $P_{T_{B}}(T)$ graph are shown.

The line degree of a line $u v$ in a tree $T$, pathos length, pathos point in $T$ was defined by Muddebihal [10]. If $G$ is planar, the inner point number $i(G)$ of a graph $G$ is the minimum number of points not belonging to the boundary of the exterior region in any embedding of $G$ in the plane. A graph $G$ is said to be minimally nonouterplanar if $i(G)=1$, as was given by Kulli [4].

We need the following results to prove further results.

Theorem [A][Ref.6] If $G$ is connected graph with $p$ points and $q$ lines and if $b_{i}$ is the number of blocks to which $v_{i}$ belongs in $G$, then the semitotal block graph $T_{b}(G)$ has $\left(\sum_{i=1}^{p} b_{i}\right)+1$, points and $q+\left(\sum_{i=1}^{p} b_{i}\right)$ lines.

Theorem [B][Ref.6] If $G$ is connected graph with $p$ points and $q$ lines and if $b_{i}$ is the number of blocks to which $v_{i}$ belongs in $G$, then the total block graph $T_{B}(G)$ has $\left(\sum_{i=1}^{p} b_{i}\right)+1$, points and $q+\sum_{i=1}^{p}\binom{b_{i}+1}{2}$ lines.

Theorem [C][Ref.8] The total block graph $T_{B}(G)$ of a graph $G$ is planar if and only if $G$ is outerplanar and every cutpoint of $G$ lies on atmost three blocks.

Theorem [D] [Ref.7] The total block graph $T_{B}(G)$ of a connected graph $G$ is minimally nonouter planar if and only if,
(1) $G$ is a cycle, or
(2) $G$ is a path $P$ of length $n \geq 2$, together with a point which is adjacent to any two adjacent points of $P$.


Figure 1:

Theorem [E][Ref.9] The total block graph $T_{B}(G)$ of a graph $G$ crossing number 1 if and only if
(1) $G$ is outer planar and every cut point in $G$ lies on at most 4 blocks and $G$ has a unique cut point which lies on 4 blocks, or
(2) $G$ is minimally non-outer planar, every cut point of $G$ lies on at most 3 blocks and exactly one block of $G$ is theta-minimally non-outer planar.

Corollary [A][Ref.1] Every nontrivial tree contains at least two end points.
Theorem [F][Ref.1] Every maximal outerplanar graph $G$ with $p$ points has $(2 p-3)$ lines.
Theorem [G][Ref.5] A graph $G$ is a non empty path if and only if it is connected graph with $p \geq 2$ points and $\sum_{i=1}^{p} d_{i}{ }^{2}-4 p+6=0$.

## §2. Pathos Semitotal Block Graph of a Tree

We start with a few preliminary results.

Remark 1 The number of blocks in pathos semitotal block graph of $P_{T_{b}}(T)$ of a tree $T$ is equal to the number of pathos in $T$.

Remark 2 If the degree of a pathos point in pathos semi total block graph $P_{T_{b}}(T)$ of a tree $T$ is $n$, then the pathos length of the corresponding path $P_{i}$ of pathos in $T$ is $n-1$.

Kulli [6] developed the new concept in graph valued functions i.e., semi total and total block graph of a graph. In this article the number of points and lines of a semi total block graph of a graph has been expressed in terms of blocks of $G$. Now using this we have a modified theorem as shown below in which we have expressed the number of points and lines in terms of lines and degrees of the points of $G$ which is a tree.

Theorem 1 For any $(p, q)$ tree $T$, the semitotal block graph $T_{b}(T)$ has $(2 q+1)$ points and $3 q$ lines.

Proof By Theorem [A], the number of points in $T_{b}(G)$ is $\left(\sum_{i=1}^{p} b_{i}\right)+1$, where $b_{i}$ are the number of blocks in $T$ to which the points $v_{i}$ belongs in $G$. Since $\sum b_{i}=2 q$, for $G$ is a tree. Thus the number of points in $T_{b}(G)=2 q+1$. Also, by Theorem [A] the number of lines in $T_{b}(G)$ are $q+\left(\sum_{i=1}^{b} b_{i}\right)$, since $\sum b_{i}=2 q$ for $G$ is a tree. Thus the number of lines in $T_{b}(G)$ is $q+2 q=3 q$.

In the following theorem we obtain the number of points and lines in $P_{T_{b}}(T)$.
Theorem 2 For any non trivial tree $T$, the pathos semitotal block graph of a tree $T$, whose points have degree $d_{i}$, then the number of points in $P_{T_{b}}(T)$ are $(2 q+k+1)$ and the number of lines are $\left(2 q+2+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$, where $k$ is the path number.

Proof By Theorem 1, the number of points in $T_{b}(T)$ are $2 q+1$, and by definition of $P_{T_{b}}(T)$, the number of points in $(2 \mathrm{q}+\mathrm{k}+1)$, where $k$ is the path number. Also by Theorem 1 , the number of lines in $T_{b}(T)$ are $3 q$. The number of lines in $P_{T_{b}}(T)$ is the sum of lines in $T_{b}(T)$ and the number of lines which lie on the points of pathos of $T$ which are to $\left(-q+2+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$. Thus the number of lines in is equal to $\left(3 q+\left(-q+2+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)\right)=\left(2 q+2+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$.

## §2. Planar Pathos Semitotal Block Graphs

A criterion for pathos semi total block graph to be planar is presented in our next theorem.
Theorem 3 For any non trivial tree $T$, the pathos semitotal block graph $P_{T_{b}}(T)$ of a tree $T$ is planar.

Proof Let $T$ be a non trivial tree, then in $T_{b}(T)$ each block is a triangle. We have the following cases.

Case 1 Suppose $G$ is a path, $G=P_{n}: u_{1}, u_{2}, u_{3}, \ldots, u_{n}, n>1$. Further, $V\left[T_{b}(T)\right]=$
$\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{n}, b_{1}, b_{2}, b_{3}, \ldots, b_{n-1}\right\}$, where $b_{1}, b_{2}, b_{3}, \ldots, b_{n-1}$ are the corresponding block points. In pathos semi total block graph $P_{T_{b}}(T)$ of a tree $T,\left\{u_{1} b_{1} u_{2} w, u_{2} b_{2} u_{3} w, u_{3} b_{3} u_{4} w, \ldots\right.$, $\left.u_{n-1} b_{n-1} u_{n} w\right\} \in V\left[P_{T_{b}}(T)\right]$, each set $\left\{u_{n-1} b_{n-1} u_{n} w\right\}$ forms an induced subgraph as $K_{4}-x$. Hence one can easily verify that $P_{T_{b}}(T)$ is planar.

Case 2 Suppose $G$ is not a path. Then $V\left[T_{b}(G)\right]=\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{n}, b_{1}, b_{2}, b_{3}, \ldots, b_{n-1}\right\}$ and $w_{1}, w_{2}, w_{3}, \ldots, w_{k}$ be the pathos points. Since $u_{n-1} u_{n}$ is a line and $u_{n-1} u_{n}=b_{n-1} \in V\left[T_{b}(G)\right]$. Then in $P_{T_{b}}(G)$ the set $\left\{u_{n-1} b_{n-1} u_{n} w\right\} \forall n>1$, forms $K_{4}-x$ as an induced subgraphs. Hence $P_{T_{b}}(G)$ is planar.

Further we develop the maximal outer planarity of $P_{T_{b}}(G)$ in the following theorem.

Theorem 4 For any non trivial tree $T$, the pathos semitotal block graph $P_{T_{b}}(T)$ of a tree $T$ is maximal outer planar if and only if $T$ is a path.

Proof Suppose $P_{T_{b}}(T)$ is maximal outer planar. Then $P_{T_{b}}(T)$ is connected. Hence $T$ is connected. If $P_{T_{b}}(T)$, is $K_{4}-x$, then obviously $T$ is $k_{2}$.

Let $T$ be any connected tree with $p \geq 2, q$ lines $b_{i}$ blocks and path number $k$, then clearly $P_{T_{b}}(T)$ has $(2 q+k+1)$ points and $\left(2 q+2+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$ lines. Since $P_{T_{b}}(T)$ is maximal outer planar, by Theorem $[\mathrm{F}]$, it has $[2(2 q+k+1)-3]$ lines. Hence,

$$
\begin{aligned}
& 2+2 q+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}=2(2 \mathrm{q}+\mathrm{k}+1)-3=4 \mathrm{q}+2 \mathrm{k}+2-3=4 \mathrm{q}+2 \mathrm{q}-1 \\
& \frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}=2 \mathrm{q}+2 \mathrm{k}-3 \\
& \sum_{i=1}^{p} d_{i}^{2}=4 \mathrm{q}+4 \mathrm{k}-6 \\
& \sum_{i=1}^{p} d_{i}^{2}=4(\mathrm{p}-1)+4 \mathrm{k}-6 \\
& \sum_{i=1}^{p} d_{i}^{2}=4 \mathrm{p}+4 \mathrm{k}-10
\end{aligned}
$$

But for a path, $k=1$.

$$
\begin{aligned}
& \sum_{i=1}^{p} d_{i}^{2}=4 \mathrm{p}+4(1)-10=4 \mathrm{p}-6 \\
& \sum_{i=1}^{p} d_{i}^{2}-4 \mathrm{p}+6=0
\end{aligned}
$$

By Theorem [G], it follows that $T$ is a non empty path. Thus necessity is proved.
For sufficiency, suppose T is a non empty path. We prove that $P_{T_{b}}(T)$ is maximal outer planar. By induction on the number of points $p_{i} \geq 2$ of $T$. It is easy to observe that $P_{T_{b}}(T)$ of a path $P$ with 2 points is $K_{4}-x$, which is maximal outer planar. As the inductive hypothesis, let the pathos semitotal block graph of a non empty path $P$ with $n$ points be maximal outer planar. We now show that the pathos semitotal block graph of a path $P^{\prime}$ with $(n+1)$ points is maximal outer planar. First we prove that it is outer planar. Let the point and line sequence of the path
$P^{\prime}$ be $v_{1}, e_{1}, v_{2}, e_{2}, v_{3}, \ldots, v_{n}, e_{n}, v_{n+1}$, Where $v_{1} v_{2}=e_{1}=b_{1}, v_{2} v_{3}=e_{2}=b_{2}, \ldots, v_{n-1} v_{n}=$ $e_{n-1}=b_{n 1}, v_{n} v_{n+1}=e_{n}=b_{n}$.

The graphs $P, P^{\prime}, T_{b}(P), T_{b}\left(P^{\prime}\right), P_{T_{b}}(P)$ and $P_{T_{b}}\left(P^{\prime}\right)$ are shown in the figure 2 . Without loss of generality $P^{\prime}-v_{n+1}=P$.

By inductive hypothesis, $P_{T_{b}}(P)$ is maximal outer planar. Now the point $v_{n+1}$ is one more point more in $P_{T_{b}}\left(P^{\prime}\right)$ than $P_{T_{b}}(P)$. Also there are only four lines $\left(v_{n+1}, v_{n}\right)\left(v_{n}, b_{n}\right)\left(b_{n}, v_{n+1}\right)$ and $\left(v_{n+1}, K_{1}\right)$ more in $P_{T_{b}}\left(P^{\prime}\right)$. Clearly the induced subgraph on the points $v_{n+1}, v_{n}, b_{n}, K_{1}$ is not $K_{4}$. Hence $P_{T_{b}}\left(P^{\prime}\right)$ is outer planar.

We now prove that $P_{T_{b}}\left(P^{\prime}\right)$ is maximal outer planar. Since $P_{T_{b}}(P)$ is maximal outer planar, it has $2(2 q+k+1)-3$ lines. The outer planar graph $P_{T_{b}}\left(P^{\prime}\right)$ has $2(2 q+k+1)-3+4=$ $2(2 q+k+1+2)-3$

$$
=2[(2 q+1)+(k+1)+1]-3 \text { lines. }
$$

By Theorem $[\mathrm{F}], P_{T_{b}}\left(P^{\prime}\right)$ is maximal outer planar.
The next theorem gives a non-minimally non-outer planar $P_{T_{b}}(T)$.
Theorem 5 For any non trivial tree $T$, the pathos semitotal block graph $P_{T_{b}}(T)$ of a tree $T$ is non-minimally non-outer planar.

Proof We have the following cases.
Case 1 Suppose $T$ is a path, then $\Delta(T) \leq 2$, then by Theorem $4, P_{T_{B}}(T)$ is maximal outer planar.

Case 2 Suppose $T$ is not a path, then $\Delta(T) \geq 3$, then by theorem $3, P_{T_{b}}(T)$ is planar. On embedding $P_{T_{b}}(T)$ in any plane, the points with degree greater than two of $T$ forms the cut points. In $P_{T_{b}}(T)$ which lie on at least two blocks. Since each block of $P_{T_{b}}(T)$ is a maximal outer planar than one can easily verify that $P_{T_{b}}(T)$ is outer planar. Hence for any non trivial tree with $\Delta(T) \geq 3, P_{T_{b}}(T)$ is non minimally non-outer planar.

In the next theorem, we characterize the non-Eulerian $P_{T_{b}}(T)$.
Theorem 6 For any non trivial tree $T$, the pathos semitotal block graph $P_{T_{b}}(T)$ of a tree $T$ is non-Eulerian.

Proof We have the following cases.
Case 1 Suppose $T$ is a path with 2 points, then $P_{T_{b}}(T)=K_{4}-x$, which is non-Eulerian. If $T$ is a path with $p>2$ points. Then in $T_{b}(T)$ each block is a triangle such that they are in sequence with the vertices of $T_{b}(T)$ as $\left\{v_{1}, b_{1}, v_{2}, v_{1}\right\}$ an induced subgraph as a triangle $T_{b}(T)$. Further $\left\{v_{2}, b_{2}, v_{3}, v_{2}\right\},\left\{v_{3}, b_{3}, v_{4}, v_{3}\right\}, \ldots,\left\{v_{n-1}, b_{n}, v_{n}, v_{n-1}\right\}$, in which each set form a triangle as an induced subgraph of $T_{b}(T)$. Clearly one can easily verify that $T_{b}(T)$ is Eulerian. Now this path has exactly one pathos point say $k_{1}$, which is adjacent to $v_{1}, v_{2}, v_{3}, \ldots, v_{n}$ in $P_{T_{b}}(T)$ in which all the points $v_{1}, v_{2}, v_{3}, \ldots, v_{n} \in P_{T_{b}}(T)$ are of odd degree. Hence $P_{T_{b}}(T)$ is non-Eulerian.

Case 2 Suppose $\Delta(T) \geq 3$. Assume $T$ has a unique point of degree $\geq 3$ and also assume that $T=K_{1 . n}$. Then in $T_{b}(T)$ each block is a triangle, such that the number of blocks which are $K_{3}$

$T_{b}(P):$

$T_{b}\left(P^{\prime}\right):$


Figure 2:
are $n$ with a common cut point $k$. Since the degree of a vertex $k=2 n$. One can easily verify that $T_{b}\left(K_{1,3}\right)$ is Eulerian. To form $P_{T_{b}}(T), T=K_{1, n}$, the points of degree 2 and the point $k$ are joined by the corresponding pathos point which give $(n+1)$ points of odd degree in $P_{T_{b}}(T)$. Hence $P_{T_{b}}(T)$ is non-Eulerian.

In the next theorem we characterize the hamiltonian $P_{T_{b}}(T)$.

Theorem 7 For any non trivial tree $T$, the pathos semitotal block graph $P_{T_{b}}(T)$ of a tree $T$ is hamiltonian if and only if $T$ is a path.

Proof For the necessity, suppose $T$ is a path and has exactly one path of pathos. Let $V\left[T_{b}(T)\right]=\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{n}\right\} \cup\left\{b_{1}, b_{2}, b_{3}, \ldots, b_{n-1}\right\}$, where $b_{1}, b_{2}, b_{3}, \ldots, b_{n-1}$ are block points of $T$. Since each block is a triangle and each block consists of points as $B_{1}=\left\{u_{1}, b_{1}, u_{2}\right\}, B_{2}=$ $\left\{u_{2}, b_{2}, u_{3}\right\}, \ldots, B_{m}=\left\{u_{m}, b_{m}, u_{m+1}\right\}$. In $P_{T_{b}}(T)$ the pathos point $w$ is adjacent to $\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{n}\right\}$. Hence $V\left[P_{T_{b}}(T)\right]=\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{n}\right\} \cup\left\{b_{1}, b_{2}, b_{3}, \ldots, b_{n-1}\right\} \cup w$ form a cycle as $w, u_{1}, b_{1}, u_{2}, b_{2}, u_{2}, \ldots$ $u_{n-1}, b_{n-1}, u_{n}, w$. Containing all the points of $P_{T_{b}}(T)$. Clearly $P_{T_{b}}(T)$ is hamiltonian. Thus necessity is proved.

For the sufficiency, suppose $P_{T_{b}}(T)$ is hamiltonian, now we consider the following cases.
Case 1 Assume $T$ is a path. Then $T$ has at least one point with $\operatorname{deg} v \geq 3, \forall v \in V(T)$, assume that $T$ has exactly one point $u$ such that degree $u>2$, then $G=T=K_{1 . n}$. Now we consider the following subcases of Case 1.

Subcase 1.1 For $K_{1 . n}, n>2$ and $n$ is even, then in $T_{b}(T)$ each block is $k_{3}$. The number of path of pathos are $\frac{n}{2}$. Since $n$ is even we get $\frac{n}{2}$ blocks. Each block contains two lines of $\left\langle K_{4}-x\right\rangle$, which is a non line disjoint subgraph of $P_{T_{b}}(T)$. Since $P_{T_{b}}(T)$ has a cut point, one can easily verify that there does not exist any hamiltonian cycle, a contradiction.

Subcase 1.2 For $K_{1 . n}, n>2$ and $n$ is odd, then the number of path of pathos are $\frac{n+1}{2}$, since n is odd we get $\frac{n-1}{2}+1$ blocks in which $\frac{n-1}{2}$ blocks contains two times of $\left\langle K_{4}-x\right\rangle$ which is nonline disjoint subgraph of $P_{T_{b}}(T)$ and remaining block is $\left\langle K_{4}-x\right\rangle$. Since $P_{T_{b}}(T)$ contain a cut point, clearly $P_{T_{b}}(T)$ does not contain a hamiltonian cycle, a contradiction. Hence the sufficient condition.

## §3. Pathos Total Block Graph of a Tree

A tree $T$, its total block graph $T_{B}(T)$, and their pathos total block graphs $P_{T_{B}}(T)$ are shown in the Fig.3. We start with a few preliminary results.

Remark 3 For any non trivial path, the inner point number of the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$ is equal to the number of cut points in $T$.

Remark 4 The degree of a pathos point in $P_{T_{B}}(T)$ is $n$, then the pathos length of the corresponding path $P_{i}$ of pathos in $T$ is $n-1$.

Remark 5 For any non trivial tree $T, P_{T_{B}}(T)$ is a block.



has been expressed in terms of blocks of $G$. Now using this we have a modified theorem as shown below in which we have expressed the number of points and lines in terms of lines and degrees of the points of $G$ which is a tree.

Theorem 8 For any non trivial $(p, q)$ tree whose points have degree $d_{i}$, the number of points and lines in total block graph of a tree $T$ are $(2 q+1)$ and $\left(2 q+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$.

Proof By Theorem [B], the number of points in $T_{b}(T)$ is $\left(\sum_{i=1}^{b} b_{i}\right)+1$, where $b_{i}$ are the number of blocks in $T$ to which the points $v_{i}$ belongs in $G$. Since $\sum b_{i}=2 q$, for $G$ is a tree. Thus the number of points in $T_{B}(G)=2 q+1$. Also, by Theorem [B], the number of lines in $T_{B}(G)$ are $q+\sum_{i=1}^{b}\binom{b_{i}+1}{2}=\left(\sum_{i=1}^{b} b_{i}\right)+\left(\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)=\left(2 q+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$, for $G$ is a tree.

In the following theorem we obtain the number of points and lines in $P_{T_{B}}(T)$.

Theorem 9 For any non trivial tree $T$, the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$, whose points have degree $d_{i}$, then the number of points in $P_{T_{B}}(T)$ are $(2 q+k+1)$ and the number of lines are $\left(q+2+\sum_{i=1}^{p} d_{i}^{2}\right)$, where $k$ is the path number.

Proof By Theorem 7, the number of points in $T_{B}(T)$ are $2 q+1$, and by definition of $P_{T_{B}}(T)$, the number of points in $P_{T_{B}}(T)$ are $(2 q+k+1)$, where $k$ is the path number in $T$. Also by Theorem 7 , the number of lines in $T_{B}(T)$ are $\left(2 q+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$. The number of lines in $P_{T_{B}}(T)$ is the sum of lines in $T_{B}(T)$ and the number of lines which lie on the points of pathos of $T$ which are to $\left(-q+2+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)$. Thus the number of lines in $P_{T_{B}}(T)$ is equal to $\left(2 q+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)+\left(-q+2+\frac{1}{2} \sum_{i=1}^{p} d_{i}^{2}\right)=\left(q+2+\sum_{i=1}^{p} d_{i}^{2}\right)$.

## §4. Planar Pathos Total Block Graphs

A criterion for pathos total block graph to be planar is presented in our next theorem.

Theorem 10 For any non trivial tree $T$, the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$ is planar if and only if $\Delta(T) \leq 3$.

Proof Suppose $P_{T_{B}}(T)$ is planar. Then by Theorem [C], each cut point of $T$ lie on at most 3 blocks. Since each block is a line in a tree, now we can consider the degree of cutpoints instead of number of blocks incident with the cut points. Now suppose if $\Delta(T) \leq 3$, then by Theorem [C], $T_{B}(T)$ is planar. Let $\left\{b_{1}, b_{2}, b_{3}, \ldots, b_{p-1}\right\}$ be the blocks of $T$ with $p$ points such that $b_{1}=e_{1}, b_{2}=e_{2}, \ldots, b_{p-1}=e_{p-1}$ and $P_{i}$ be the number of pathos of $T$. Now $V\left[P_{T_{B}}(T)\right]=V(G) \cup\left\{b_{1}, b_{2}, \ldots b_{p-1}\right\} \cup\left\{P_{i}\right\}$. By Theorem [C], and by the definition of pathos, the embedding of $P_{T_{B}}(T)$ in any plane gives a planar $P_{T_{B}}(T)$.

Suppose $\Delta(T) \geq 4$ and assume that $P_{T_{B}}(T)$ is planar. Then there exists at least one point
of degree 4, assume that there exists a vertex $v$ such that $\operatorname{deg} v=4$. Then in $T_{B}(T)$, this point together with the block points form $k_{5}$ as an induced subgraph. Further the corresponding pathos point are adjacent to the $\mathrm{V}(\mathrm{T})$ in $T_{B}(T)$ which gives $P_{T_{B}}(T)$ in which again $k_{5}$ as an induced subgraph, a contradiction to the planarity of $P_{T_{B}}(T)$. This completes the proof.

The following theorem results the maximal outer planar $P_{T_{B}}(T)$.
Theorem 11 For any non trivial tree $T$, the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$ is maximal outer planar if and only if $T=k_{2}$.

Proof Suppose $T=k_{3}$ and $P_{T_{B}}(T)$ is maximal outer planar. Then $T_{B}(T)=k_{4}$ and one can easily verify that, $i\left[P_{T_{B}}(T)\right]>1$, a contradiction. Further we assume that $T=K_{1,2}$ and $P_{T_{B}}(T)$ is maximal outer planar, then $T_{B}(T)$ is $W_{p}-x$, where $x$ is outer line of $W_{p}$. Since $K_{1,2}$ has exactly one pathos, this point together with $W_{p}-x$ gives $W_{p+1}$. Clearly, $P_{T_{B}}(T)$ is non maximal outer planar, a contradiction. For the converse, if $T=k_{2}, T_{B}(T)=k_{3}$ and $P_{T_{B}}(T)$ $=K_{4}-x$ which is a maximal outer planar. This completes the proof of the theorem.

Now we have a pathos total block graph of a path $p \geq 2$ point as shown in the below remarks, and also a cycle with $p \geq 3$ points.

Remark 6 For any non trivial path with $p$ points, $i\left[P_{T_{B}}(T)\right]=p-2$.
Remark 7 For any cycle $C_{p}, p \geq 3, i\left[P_{T_{B}}\left(C_{p}\right)\right]=p-1$.
The next theorem gives a minimally non-outer planar $P_{T_{B}}(T)$.

Theorem 12 For any non trivial tree $T$, the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$ is minimally non-outer planar if and only if $T$ is a path with 3 points.

Proof Suppose $P_{T_{B}}(T)$ is minimally non-outer planar. Assume $T$ is not a path. We consider the following cases.

Case 1 Suppose $T$ is a tree with $\Delta(T) \geq 3$. Then there exists at least one point of degree at least 3. Assume $v$ be a point of degree 3. Clearly, $T=K_{1,3}$. Then by the Theorem [D], $i\left[T_{B}(T)\right]>1$ since $T_{B}(T)$ is a subgraph of $P_{T_{B}}(T)$. Clearly $i\left[P_{T_{B}}(T)\right] \geqslant 2$ a contradiction.

Case 2 Suppose $T$ is a closed path with $p$ points, then it is a cycle with $p$ points. By Theorem $[\mathrm{D}], P_{T_{B}}(T)$ is minimally non-outer planar. By Remark $7, i\left[P_{T_{B}}(T)\right]>1$, a contradiction.

Case 3 Suppose $T$ is a closed path with $p \geq 4$ points, clearly by Remark $6, i\left[P_{T_{B}}(T)\right]>2$, a contradiction.

Conversely, suppose $T$ is a path with 3 points, clearly by Remark $6, i\left[P_{T_{B}}(T)\right]=1$. This gives the required result.

In the following theorem we characterize the non-Eulerian $P_{T_{B}}(T)$.
Theorem 13 For any non trivial tree $T$, the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$ is non-Eulerian.

Proof We consider the following cases.
Case 1 Suppose $T$ is a path. For $p=2$ points, then $P_{T_{B}}(T)=K_{4}-x$, which is non-Eulerian. For $p=3$ points, then $P_{T_{B}}(T)$ is a wheel, which is non-Eulerian.

For $p \geq 4$ we have a path $P: v_{1}, v_{2}, v_{3}, \ldots, v_{p}$. Now in path each line is a block. Then $v_{1} v_{2}=e_{1}=b_{1}, v_{2} v_{3}=e_{2}=b_{2}, \ldots, v_{p-1} v_{p}=e_{p-1}=b_{p-1}, \forall e_{p-1} \in E(G)$, and $\forall b_{p-1} \in$ $V\left[T_{B}(P)\right]$. In $T_{B}(P)$, the degree of each point is even except $b_{1}$ and $b_{p-1}$. Since the path $P$ has exactly one pathos which is a point of $P_{T_{B}}(P)$ and is adjacent to the points $v_{1}, v_{2}, v_{3}, \ldots, v_{p}$, of $T_{B}(P)$ which are of even degree, gives as an odd degree points in $P_{T_{B}}(P)$ including odd degree points $b_{1}$ and $b_{2}$. Clearly $P_{T_{B}}(P)$ is non-Eulerian.

Case 2 Suppose $T$ is not a path. We consider the following subcases,
Subcase 2.1 Assume $T$ has a unique point degree $\geq 3$ and $T=K_{1 . n}$, where $n$ is odd. Then in $T_{B}(T)$ each block is a triangle such that there are $n$ number of triangles with a common cut point $k$ which has a degree $2 n$. Since the degree of each point in $T_{B}\left(K_{1, n}\right)$ is Eulerian. To form $P_{T_{B}}(T)$ where $T=K_{1, n}$, the points of degree 2 and the point $k$ are joined by the corresponding pathos point which gives $(n+1)$ points of odd degree in $P_{T_{B}}\left(K_{1 . n}\right) . P_{T_{B}}\left(K_{1 . n}\right)$ has $n$ points of odd degree. Hence $P_{T_{B}}(T)$ non-Eulerian.

Assume that $T=K_{1 . n}$, where $n$ is even, then in $T_{B}(T)$ each block is a triangle, which are $2 n$ in number with a common cut point $k$. Since the degree of each point other than $k$ is either 2 or $(n+1)$ and the degree of the point $k$ is $2 n$. One can easily verify that $T_{B}\left(K_{1, n}\right)$ is non-Eulerian. To form $P_{T_{B}}(T)$ where $T=K_{1, n}$, the points of degree 2 and the point $k$ are joined by the corresponding pathos point which gives $(n+2)$ points of odd degree in $P_{T_{B}}(T)$. Hence $P_{T_{B}}(T)$ non-Eulerian.

Subcase 2.2 Assume $T$ has at least two points of degree $\geq 3$. Then $V\left[T_{B}(T)\right]=V(G) \cup$ $\left\{b_{1}, b_{2}, b_{3}, \ldots, b_{p}\right\}, \forall e_{p} \in E(G)$. In $T_{B}(T)$, each endpoint has degree 2 and these points are adjacent to the corresponding pathos points in $P_{T_{B}}(T)$ gives degree 3, From Case 1, Tree $T$ has at least 4 points and by Corollary [A], $P_{T_{B}}(T)$ has at least two points of degree 3 . Hence $P_{T_{B}}(T)$ is non-Eulerian.

In the next theorem we characterize the hamiltonian $P_{T_{B}}(T)$.
Theorem 14 For any non trivial tree $T$, the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$ is hamiltonian.

Proof We consider the following cases.
Case 1 Suppose $T$ is a path with $\left\{u_{1}, u_{2}, u_{3}, \ldots, u_{n}\right\} \in V(T)$ and $b_{1}, b_{2}, b_{3}, \ldots, b_{m}$ be the number of blocks of $T$ such that $m=n-1$. Then it has exactly one path of pathos. Now point set of $T_{B}(T)$ is $V\left[T_{B}(T)\right]=\left\{u_{1}, u_{2}, \ldots, u_{n}\right\} \cup\left\{b_{1}, b_{2}, \ldots, b_{m}\right\}$. Since given graph is a path then in $T_{B}(T), b_{1}=e_{1}, b_{2}=e_{2}, \ldots, b_{m}=e_{m}$, such that $b_{1}, b_{2}, b_{3}, \ldots, b_{m} \subset V\left[T_{B}(T)\right]$. Then by the definition of total block graph $\left\{u_{1}, u_{2}, \ldots, u_{m}\right\} \cup\left\{b_{1}, b_{2}, \ldots, b_{m-1}, b_{m}\right\} \cup\left\{b_{1}, u_{1}, b_{2} u_{2}, \ldots, b_{m} u_{n-1}\right.$, $\left.b_{m} u_{n}\right\}$ form line set of $T_{B}(T)$ [see Fig. 4].

Now this path has exactly one pathos say $w$. In forming pathos total block graph of a path, the pathos $w$ becomes a point, then $V\left[P_{T_{B}}(T)\right]=\left\{u_{1}, u_{2}, \ldots, u_{n}\right\} \cup\left\{b_{1}, b_{2}, \ldots, b_{m}\right\} \cup\{w\}$ and


Figure 4:
$w$ is adjacent to all the points $\left\{u_{1}, u_{2}, \ldots, u_{m}\right\}$ shown in the Fig.5.


Figure 5:

In $P_{T_{B}}(T)$, the hamiltonian cycle $w, u_{1}, b_{1}, u_{2}, b_{2}, u_{2}, u_{3}, b_{3}, \ldots, u_{n-1}, b_{m}, u_{n}, w$ exist. Clearly the pathos total block graph of a path is hamiltonian graph.

Case 2 Suppose $T$ is not a path. Then $T$ has at least one point with degree at least 3. Assume that $T$ has exactly one point $u$ such that degree $>2$. Now we consider the following subcases of case 2 .

Subcase 2.1 Assume $T=K_{1 . n}, n>2$ and is odd. Then the number of paths of pathos are $\frac{n+1}{2}$. Let $V\left[T_{B}(T)\right]=\left\{u_{1}, u_{2}, \ldots, u_{n}, b_{1}, b_{2}, \ldots, b_{m-1}\right\}$. By the definition of $P_{T_{B}}(T)$, $V\left[P_{T_{B}}(T)\right]=\left\{u_{1}, u_{2}, \ldots, u_{n}, b_{1} b_{2}, \ldots, b_{n-1}\right\} \cup\left\{p_{1}, p_{2}, \ldots, p_{n+1 / 2}\right\}$. Then there exists a cycle containing the points of $P_{T_{B}}(T)$ as $p_{1}, u_{1}, b_{1}, b_{2}, u_{3}, p_{2}, u_{2}, b_{3}, u_{4}, \ldots p_{1}$ and is a hamiltonian cycle. Hence $P_{T_{B}}(T)$ is a hamiltonian.

Subcase 2.2 Assume $T=K_{1 . n}, n>2$ and is even. Then the number of path of pathos are $\frac{n}{2}$, then $V\left[T_{B}(T)\right]=\left\{u_{1}, u_{2}, \ldots, u_{n}, b_{1}, b_{2}, \ldots b_{n-1}\right\}$. By the definition of $P_{T_{B}}(T) . V\left[P_{T_{B}}(T)\right]=$ $\left\{u_{1}, u_{2}, \ldots, u_{n}, b_{1}, b_{2}, \ldots, b_{n-1}\right\} \cup\left\{p_{1}, p_{2}, \ldots, p_{n / 2}\right\}$. Then there exist a cycle containing the points of $P_{T_{B}}(T)$ as $p_{1}, u_{1}, b_{1}, b_{2}, u_{3}, p_{2}, u_{4}, b_{3}, b_{4}, \ldots, p_{1}$ and is a hamiltonian cycle. Hence $P_{T_{B}}(T)$ is a hamiltonian.

Suppose $T$ is neither a path or a star. Then $T$ contains at least two points of degree $>2$. Let $u_{1}, u_{2}, u_{3}, \ldots, u_{n}$ be the points of degree $\geq 2$ and $v_{1}, v_{2}, v_{3}, \ldots, v_{m}$ be the end points of $T$. Since end block is a line in $T$, and denoted as $b_{1}, b_{2}, \ldots, b_{k}$, then tree $T$ has $p_{i}$ pathos points, $i>1$ and each pathos point is adjacent to the point of $T$ where the corresponding pathos lie on the points of $T$. Let $\left\{p_{1}, p_{2}, \ldots ., p_{i}\right\}$ be the pathos points of $T$. Then there exists a cycle $C$ containing all the points of $P_{T_{B}}(T)$ as $p_{1}, v_{1}, b_{1}, b_{2}, v_{2}, p_{2}, u_{1}, b_{3}, u_{2}, p_{3}, v_{3}, b_{4}, \ldots, v_{n-1}, b_{n-1}, b_{n}, v_{n}, \ldots, p_{1}$. Hence $P_{T_{B}}(T)$ is a hamiltonian cycle. Hence $P_{T_{B}}(T)$ is a hamiltonian graph.

In the next theorem we characterize $P_{T_{B}}(T)$ in terms of crossing number one.
Theorem 15 For any non trivial tree $T$, the pathos total block graph $P_{T_{B}}(T)$ of a tree $T$ has crossing number one if and only if $\Delta(T) \leq 4$, and there exist a unique point in $T$ of degree 4.

Proof Suppose $P_{T_{B}}(T)$ has crossing number one. Then it is non-planar. Then by Theorem 10, we have $\Delta(T) \geq 4$. We now consider the following cases.

Case 1 Assume $\Delta(T)=5$. Then by Theorem [E], $T_{B}(T)$ is non-planar with crossing number more than one. Since $T_{B}(T)$ is a subgraph of $P_{T_{B}}(T)$. Clearly $\operatorname{cr}\left(P_{T_{B}}(T)\right)>1$, a contradiction.

Case 2 Assume $\Delta(T)=4$. Suppose $T$ has two points of degree 4. Then by Theorem [E], $T_{B}(T)$ has crossing number at least two. But $T_{B}(T)$ is a subgraph of $P_{T_{B}}(T)$. Hence $\operatorname{cr}\left(P_{T_{B}}(T)\right)>1$, a contradiction.

Conversely, suppose $T$ satisfies the given condition and assume $T$ has a unique point $v$ of degree 4. The lines which are blocks in $T$ such that they are the points in $T_{B}(T)$. In $T_{B}(T)$, these block points and a point $v$ together forms an induced subgraph as $k_{5}$. In forming $P_{T_{B}}(T)$, the pathos points are adjacent to at most two points of this induced subgraph. Hence in all these cases the $\operatorname{cr}\left(P_{T_{B}}(T)\right)=1$. This completes the proof.

## References

[1] Harary F., Graph Theory, Addison-Wesley Reading. Mass. (1969), pp 34 and 107.
[2] Harary F., Covering and packing in graphs I, Annals of New York Academy of Science, (1970), pp 198-205.
[3] Harary F. R. and Schwenk A. J., Evolution of the path number of a graph, covering and packing in graphs II, Graph Theory and Computing, Ed. R. C. Read, Academic Press, New York (1972), pp 39-45.
[4] Kulli V. R., On minimally non-outer planar graphs, Proceeding of the Indian National Science Academy, Vol.41, Part A, No.3, (1975), pp 275-280.
[5] Kulli V. R. A Characterization of Paths, Vol.9, No.1, The Maths Education, (1975), pp 1-2.
[6] Kulli V. R., The semitotal-block graph and total-block graph of a graph, Indian J. Pure and App. Math., Vol.7, No. 6 (1976), pp 625-630.
[7] Kulli V. R. and Patil H. P., Minimally non-outer planar graphs and some graph valued functions, Journal of Karnataka University Science, (1976), pp 123-129.
[8] Kulli V. R. and Akka D. R., Transversability and planarity of total block graphs, J. Math.

Phy. Sci., (1977) Vol.11, pp 365-375.
[9] Kulli V. R. and Muddebihal M. H., Total block graphs and semitotal block graphs with crossing numbers, Far East J.Appl. Math., 4(1), (2000), pp 99-106.
[10] Muddebihal M. H., Gudagudi B. R. and Chandrasekhar R., On pathos line graph of a tree, National Academy Science Letters, Vol.24, No. 5 to 12(2001), pp 116-123.
[11] Stanton R. G., Cowan D. D. and James L. O., Some results on the path numbers, Proceedings of Lousiana Conference on Combinatorics, Graph Theory and Computing, (1970), pp 112- 135.

# Varieties of Groupoids and Quasigroups 

# Generated by Linear-Bivariate Polynomials Over Ring $\mathbb{Z}_{n}$ 

E.Ilojide, T.G.Jaiyéọlá and O.O.Owojori<br>(Department of Mathematics, Obafemi Awolowo University, Ile-Ife 220005, Nigeria)

E-mail: emmailojide@yahoo.com, jaiyeolatemitope@yahoo.com, tjayeola@oauife.edu.ng, walejori@yahoo.com


#### Abstract

Some varieties of groupoids and quasigroups generated by linear-bivariate polynomials $P(x, y)=a+b x+c y$ over the ring $\mathbb{Z}_{n}$ are studied. Necessary and sufficient conditions for such groupoids and quasigroups to obey identities which involve one, two, three (e.g. BolMoufang type) and four variables w.r.t. $a, b$ and $c$ are established. Necessary and sufficient conditions for such groupoids and quasigroups to obey some inverse properties w.r.t. $a, b$ and $c$ are also established. This class of groupoids and quasigroups are found to belong to some varieties of groupoids and quasigroups such as medial groupoid(quasigroup), F-quasigroup, semi automorphic inverse property groupoid(quasigroup) and automorphic inverse property groupoid(quasigroup).


Key Words: groupoids, quasigroups, linear-bivariate polynomials.
AMS(2010): 20N02, 20NO5

## §1. Introduction

### 1.1 Groupoids, Quasigroups and Identities

Let $G$ be a non-empty set. Define a binary operation $(\cdot)$ on $G .(G, \cdot)$ is called a groupoid if $G$ is closed under the binary operation $(\cdot)$. A groupoid $(G, \cdot)$ is called a quasigroup if the equations $a \cdot x=b$ and $y \cdot c=d$ have unique solutions for $x$ and $y$ for all $a, b, c, d \in G$. A quasigroup $(G, \cdot)$ is called a loop if there exists a unique element $e \in G$ called the identity element such that $x \cdot e=e \cdot x=x$ for all $x \in G$.

A function $f: S \times S \rightarrow S$ on a finite set $S$ of size $n>0$ is said to be a Latin square (of order $n$ ) if for any value $a \in S$ both functions $f(a, \cdot)$ and $f(\cdot, a)$ are permutations of $S$. That is, a Latin square is a square matrix with $n^{2}$ entries of $n$ different elements, none of them occurring more than once within any row or column of the matrix.

Definition 1.1 A pair of Latin squares $f_{1}(\cdot, \cdot)$ and $f_{2}(\cdot, \cdot)$ is said to be orthogonal if the pairs $\left(f_{1}(x, y), f_{2}(x, y)\right)$ are all distinct, as $x$ and $y$ vary.

For associative binary systems, the concept of an inverse element is only meaningful if the

[^7]system has an identity element. For example, in a group ( $G, \cdot$ ) with identity element $e \in G$, if $x \in G$ then the inverse element for $x$ is the element $x^{-1} \in G$ such that
$$
x \cdot x^{-1}=x^{-1} \cdot x=e
$$

In a loop $(G, \cdot)$ with identity element $e$, the left inverse element of $x \in G$ is the element $x^{\lambda} \in G$ such that

$$
x^{\lambda} \cdot x=e
$$

while the right inverse element of $x \in G$ is the element $x^{\rho} \in G$ such that

$$
x \cdot x^{\rho}=e
$$

In case $(G, \cdot)$ is a quasigroup, then for each $x \in G$, the elements $x^{\rho} \in G$ and $x^{\lambda} \in G$ such that $x x^{\rho}=e^{\rho}$ and $x^{\lambda} x=e^{\lambda}$ are called the right and left inverse elements of $x$ respectively. Here, $e^{\rho} \in G$ and $e^{\lambda} \in G$ satisfy the relations $x e^{\rho}=x$ and $e^{\lambda} x=x$ for all $x \in G$ and are respectively called the right and left identity elements. Whenever $e^{\rho}=e^{\lambda}$, then $(G, \cdot)$ becomes a loop.

In case $(G, \cdot)$ is a groupoid, then for each $x \in G$, the elements $x^{\rho} \in G$ and $x^{\lambda} \in G$ such that $x x^{\rho}=e_{\rho}(x)$ and $x^{\lambda} x=e_{\lambda}(x)$ are called the right and left inverse elements of $x$ respectively. Here, $e_{\rho}(x) \in G$ and $e_{\lambda}(x) \in G$ satisfy the relations $x e_{\rho}(x)=x$ and $e_{\lambda}(x) x=x$ for each $x \in G$ and are respectively called the local right and local left identity elements of $x$. Whenever $e_{\rho}(x)=e_{\lambda}(x)$, then we simply write $e(x)=e_{\rho}(x)=e_{\lambda}(x)$ and call it the local identity of $x$.

The basic text books on quasigroups, loops are Pflugfelder [19], Bruck [1], Chein, Pflugfelder and Smith [2], Dene and Keedwell [3], Goodaire, Jespers and Milies [4], Sabinin [25], Smith [26], Jaíyéọlá [5] and Vasantha Kandasamy [28].

Groupoids, quasigroups and loops are usually studied relative to properties or identities. If a groupoid, quasigroup or loop obeys a particular identity, then such types of groupoids, quasigroups or loops are said to form a variety. In this work, our focus will be on groupoids and quasigroups. Some identities that describe groupoids and quasigroups which would be of interest to us here are categorized as follows:
(A) Those identities which involve one element only on each side of the equality sign:

$$
\begin{array}{lc}
a a=a & \text { idempotent law } \\
a a=b b & \text { unipotent law } \tag{2}
\end{array}
$$

(B) Those identities which involve two elements on one or both sides of the equality sign:

$$
\begin{array}{cc}
a b=b a & \text { commutative law } \\
(a b) b=a & \text { Sade right Keys law } \\
b(b a)=a & \text { Sade left keys law } \\
(a b) b=a(b b) & \text { right alternative law } \\
b(b a)=(b b) a & \text { left alternative law } \tag{7}
\end{array}
$$

$$
\begin{array}{cc}
a(b a)=(a b) a & \text { medial alternative law } \\
a(b a)=b & \text { law of right semisymmetry } \\
(a b) a=b & \text { law of left semisymmetry } \\
a(a b)=b a & \text { Stein first law } \\
a(b a)=(b a) a & \text { Stein second law } \\
a(a b)=(a b) b & \text { Schroder first law } \\
(a b)(b a)=a & \text { Schroder second law } \\
(a b)(b a)=b & \text { Stein third law } \\
a b=a & \text { Sade right translation law } \\
a b=b & \text { Sade left translation law } \tag{17}
\end{array}
$$

(C) Those identities which involve three distinct elements on one or both sides of the equality sign:

$$
\begin{array}{cc}
(a b) c=a(b c) & \text { associative law } \\
a(b c)=c(a b) & \text { law of cyclic associativity } \\
(a b) c=(a c) b & \text { law of right permutability } \\
a(b c)=b(a c) & \text { law of left permutability } \\
a(b c)=c(b a) & \text { Abel-Grassman law } \\
(a b) c=a(c b) & \text { commuting product law } \\
c(b a)=(b c) a & \text { dual of commuting product } \\
(a b)(b c)=a c & \text { Stein fourth law } \\
(b a)(c a)=b c & \text { law of right transitivity } \\
(a b)(a c)=b c & \text { law of left transitivity } \\
(a b)(a c)=c b & \text { Schweitzer law } \\
(b a)(c a)=c b & \text { dual of Schweitzer law } \\
(a b) c=(a c)(b c) & \text { law of right self-distributivity law } \\
c(b a)=(c b)(c a) \\
(a b) c=(c a)(b c) & \text { law of left self-distributivity law } \\
c(b a)=(c b)(a c) & \text { law of right abelian distributivity } \\
(a b)(c a)=[a(b c)] a & \text { law of left abelian distributivity } \\
(a b)(c a)=a[b c) a] & \text { dual of Bruck-Moufang identity } \\
{[(a b) c] b=a[b(c b)]} & \text { Moufang identity } \tag{36}
\end{array}
$$

$$
\begin{align*}
{[(b c) b] a=b[c(b a)] } & \text { Moufang identity }  \tag{37}\\
{[(a b) c] b=a[(b c) b] } & \text { right Bol identity }  \tag{38}\\
{[b(c b)] a=b[c(b a)] } & \text { left Bol identity }  \tag{39}\\
{[(a b) c] a=a[b(c a)] } & \text { extra law }  \tag{40}\\
{[(b a) a] c=b[(a a) c] } & \mathrm{RC}_{4} \text { law }  \tag{41}\\
{[b(a a)] c=b[a(a c)] } & \mathrm{LC}_{4} \text { law }  \tag{42}\\
(a a)(b c)=[a(a b)] c & \mathrm{LC}_{2} \text { law }  \tag{43}\\
{[(b c) a] a=b[(c a) a] } & \mathrm{RC}_{1} \text { law }  \tag{44}\\
{[a(a b)] c=a[a(b c)] } & \mathrm{LC}_{1} \text { law }  \tag{45}\\
(b c)(a a)=b[(c a) a] & \mathrm{RC}_{2} \text { law }  \tag{46}\\
{[(a a) b] c=a[a(b c)] } & \mathrm{LC}_{3} \text { law }  \tag{47}\\
{[(b c) a] a=b[c(a a)] } & \mathrm{RC}_{3} \text { law }  \tag{48}\\
{[(b a) a] c=b[a(a c)] } & \mathrm{C}^{2}-\mathrm{law}  \tag{49}\\
a[b(c a)]=c b & \text { Tarski law }  \tag{50}\\
a[(b c)(b a)]=c & \text { Neumann law }  \tag{51}\\
(a b)(c a)=(a c)(b a) & \text { specialized medial law } \tag{52}
\end{align*}
$$

(D) Those involving four elements:

$$
\begin{array}{rc}
(a b)(c d)=(a d)(c b) & \text { first rectangle rule } \\
(a b)(a c)=(d b)(d c) & \text { second rectangle rule } \\
(a b)(c d)=(a c)(b d) & \text { internal mediality or medial law } \tag{55}
\end{array}
$$

(E) Those involving left or right inverse elements:

$$
\begin{array}{cc}
x^{\lambda} \cdot x y=y & \text { left inverse property } \\
y x \cdot x^{\rho}=y & \text { right inverse property } \\
x(y x)^{\rho}=y^{\rho} \text { or }(x y)^{\lambda} x=y^{\lambda} & \text { weak inverse property(WIP) } \tag{58}
\end{array}
$$

$x y \cdot x^{\rho}=y$ or $x \cdot y x^{\rho}=y$ or $x^{\lambda} \cdot(y x)=y$ or $x^{\lambda} y \cdot x=y$ cross inverse property(CIP) (59)
$(x y)^{\rho}=x^{\rho} y^{\rho}$ or $(x y)^{\lambda}=x^{\lambda} y^{\lambda}$ automorphic inverse property (AIP)
$(x y)^{\rho}=y^{\rho} x^{\rho}$ or $(x y)^{\lambda}=y^{\lambda} x^{\lambda}$ anti-automorphic inverse property (AAIP)
$(x y \cdot x)^{\rho}=x^{\rho} y^{\rho} \cdot x^{\rho}$ or $(x y \cdot x)^{\lambda}=x^{\lambda} y^{\lambda} \cdot x^{\lambda}$ semi-automorphic inverse property (SAIP)

Definition 1.2(Trimedial Quasigroup) A quasigroup is trimedial if every subquasigroup generated by three elements is medial.

Medial quasigroups have also been called abelian, entropic, and other names, while trimedial quasigroups have also been called triabelian, terentropic, etc.

There are two distinct, but related, generalizations of trimedial quasigroups. The variety of semimedial quasigroups(also known as weakly abelian, weakly medial, etc.) is defined by the equations

$$
\begin{align*}
& x x \cdot y z=x y \cdot x z  \tag{63}\\
& z y \cdot x x=z x \cdot y z \tag{64}
\end{align*}
$$

Definition 1.3(Semimedial Quasigroup) A quasigroup satisfying (63) (resp. (64) is said to be left (resp. right) semimedial.

Definition 1.4(Medial-Like Identities) A groupoid or quasigroup is called an external medial groupoid or quasigroup if it obeys the identity

$$
\begin{equation*}
a b \cdot c d=d b \cdot c a \quad \text { external medial or paramediality law } \tag{65}
\end{equation*}
$$

A groupoid or quasigroup is called a palindromic groupoid or quasigroup if it obeys the identity

$$
\begin{equation*}
a b \cdot c d=d c \cdot b a \quad \text { palidromity law } \tag{66}
\end{equation*}
$$

Other medial like identities of the form $(a b)(c d)=(\pi(a) \pi(b))(\pi(c) \pi(d))$, where $\pi$ is a certain permutation on $\{a, b, c, d\}$ are given as follows:

$$
\begin{array}{rlrl}
a b \cdot c d & =a b \cdot d c & & C_{1} \\
a b \cdot c d & =b a \cdot c d & & C_{2} \\
a b \cdot c d & =b a \cdot d c & & C_{3} \\
a b \cdot c d & =c d \cdot a b & & C_{4} \\
a b \cdot c d & =c d \cdot b a & C \\
a b \cdot c d & =d c \cdot a b & C & C_{6} \\
a b \cdot c d & =a c \cdot d b & C M_{1} \\
a b \cdot c d & =a d \cdot b c & C M_{2} \\
a b \cdot c d & =a d \cdot c b & C M_{3} \\
a b \cdot c d & =b c \cdot a d & C M_{4} \\
a b \cdot c d & =b c \cdot d a & C M_{5} \\
a b \cdot c d & =b d \cdot a c & C M_{6} \\
a b \cdot c d & =b d \cdot c a & C M_{7} \tag{79}
\end{array}
$$

$$
\begin{array}{rlrl}
a b \cdot c d & =c a \cdot b d & & C M_{8} \\
a b \cdot c d & =c a \cdot d b & C M_{9} \\
a b \cdot c d & =c b \cdot a d & C M_{10} \\
a b \cdot c d & =c b \cdot d a & C M_{11} \\
a b \cdot c d & =d a \cdot b c & C M_{12} \\
a b \cdot c d & =d a \cdot c b & C M_{13} \\
a b \cdot c d & =d b \cdot a c & C M_{14} \tag{86}
\end{array}
$$

The variety of F-quasigroups was introduced by Murdoch [18].

Definition 1.5(F-quasigroup) An F-quasigroup is a quasigroup that obeys the identities

$$
\begin{array}{ll}
x \cdot y z=x y \cdot(x \backslash x) z & \text { left F-law } \\
z y \cdot x=z(x / x) \cdot y x &  \tag{88}\\
\text { right F-law }
\end{array}
$$

A quasigroup satisfying (87) (resp. (77)) is called a left (resp. right) F-quasigroup.
Definition 1.6(E-quasigroup) An E-quasigroup is a quasigroup that obeys the identities

$$
\begin{array}{ll}
x \cdot y z=e_{\lambda}(x) y \cdot x z & E_{l} l a w \\
z y \cdot x=z x \cdot y e_{\rho}(x) & E_{r} L a w \tag{90}
\end{array}
$$

A quasigroup satisfying (89) (resp. (90)) is called a left (resp. right) E-quasigroup.
Some identities will make a quasigroup to be a loop, such are discussed in Keedwell [6]-[7].
Definition 1.7(Linear Quasigroup and T-quasigroup) A quasigroup $(Q, \cdot)$ of the form $x \cdot y=$ $x \alpha+y \beta+c$ where $(Q,+)$ is a group, $\alpha$ is its automorphism and $\beta$ is a permutation of the set $Q$, is called a left linear quasigroup.

A quasigroup $(Q, \cdot)$ of the form $x \cdot y=x \alpha+y \beta+c$ where $(Q,+)$ is a group, $\beta$ is its automorphism and $\alpha$ is a permutation of the set $Q$, is called a right linear quasigroup.

A T-quasigroup is a quasigroup $(Q, \cdot)$ defined over an abelian group $(Q,+)$ by $x \cdot y=$ $c+x \alpha+y \beta$, where $c$ is a fixed element of $Q$ and $\alpha$ and $\beta$ are both automorphisms of the group $(Q,+)$.

Whenever one considers mathematical objects defined in some abstract manner, it is usually desirable to determine that such objects exist. Although occasionally this is accomplished by means of an abstract existential argument, most frequently, it is carried out through the presentation of a suitable example, often one which has been specifically constructed for the purpose. An example is the solution to the open problem of the axiomization of rectangular quasigroups and loops by Kinyon and Phillips [12] and the axiomization of trimedial quasigroups by Kinyon and Phillips [10], [11].

Chein et. al. [2] presents a survey of various methods of construction which has been used in the literature to generate examples of groupoids and quasigroups. Many of these constructions are ad hoc-designed specifically to produce a particular example; while others are of more general applicability. More can be found on the construction of $(r, s, t)$-inverse quasigroups in Keedwell and Shcherbacov [8]-[9], idempotent medial quasigroups in Krc̆adinac and Volenec [14] and quasigroups of Bol-Moufang type in Kunen [15]-[16].

Remark 1.1 In the survey of methods of construction of varieties and types of quasigroups highlighted in Chein et. al. [2], it will be observed that some other important types of quasigroups that obey identities (1) to (90) are not mentioned. Also, examples of methods of construction of such varieties that are groupoids are also scarce or probably not in existence by our search. In Theorem 1.4 of Kirnasovsky [13], the author characterized T-quasigroups with a score and two identities from among identities (1) to (90). The present work thus proves some results with which such groupoids and quasigroups can be constructed.

### 1.2 Univariate and Bivariate Polynomials

Consider the following definitions.

Definition 1.8 A polynomial $P(x)=a_{0}+a_{1} x+\cdots+a_{n} x^{n}, n \in \mathbb{N}$ is said to be a permutation polynomial over a finite ring $R$ if the mapping defined by $P$ is a bijection on $R$.

Definition 1.9 A bivariate polynomial is a polynomial in two variables, $x$ and $y$ of the form $P(x, y)=\Sigma_{i, j} a_{i j} x^{i} y^{j}$.

Definition 1.10(Bivariate Polynomial Representing a Latin Square) A bivariate polynomial $P(x, y)$ over $\mathbb{Z}_{n}$ is said to represent (or generate) a Latin square if $\left(\mathbb{Z}_{n}, *\right)$ is a quasigroup where $*: \mathbb{Z}_{n} \times \mathbb{Z}_{n} \rightarrow \mathbb{Z}_{n}$ is defined by $x * y=P(x, y)$ for all $x, y \in \mathbb{Z}_{n}$.

Mollin and Small [17] considered the problem of characterizing permutation polynomials. They established conditions on the coefficients of a polynomial which are necessary and sufficient for it to represent a permutation.

Shortly after, Rudolf and Mullen [23] provided a brief survey of the main known classes of permutation polynomials over a finite field and discussed some problems concerning permutation polynomials (PPs). They described several applications of permutations which indicated why the study of permutations is of interest. Permutations of finite fields have become of considerable interest in the construction of cryptographic systems for the secure transmission of data. Thereafter, the same authors in their paper [24], described some results that had appeared after their earlier work including two major breakthroughs.

Rivest [22] studied permutation polynomials over the ring $\left(\mathbb{Z}_{n},+, \cdot\right)$ where $n$ is a power of 2: $n=2^{w}$. This is based on the fact that modern computers perform computations modulo $2^{w}$ efficiently (where $w=2,8,16,32$ or 64 is the word size of the machine), and so it was of interest to study PPs modulo a power of 2 . Below is an important result from his work which is relevant to the present study.

Theorem 1.1(Rivest [22]) A bivariate polynomial $P(x, y)=\Sigma_{i, j} a_{i j} x^{i} y^{j}$ represents a Latin square modulo $n=2^{w}$, where $w \geq 2$, if and only if the four univariate polynomials $P(x, 0)$, $P(x, 1), P(0, y)$, and $P(1, y)$ are all permutation polynomial modulo $n$.

Vadiraja and Shankar [27] motivated by the work of Rivest continued the study of permutation polynomials over the ring $\left(\mathbb{Z}_{n},+, \cdot\right)$ by studying Latin squares represented by linear and quadratic bivariate polynomials over $\mathbb{Z}_{n}$ when $n \neq 2^{w}$ with the characterization of some PPs. Some of the main results they got are stated below.

Theorem 1.2(Vadiraja and Shankar [27]) A bivariate linear polynomial $a+b x+c y$ represents a Latin square over $\mathbb{Z}_{n}, n \neq 2^{w}$ if and only if one of the following equivalent conditions is satisfied:
(i) both $b$ and $c$ are coprime with $n$;
(ii) $a+b x, a+c y,(a+c)+b x$ and $(a+b)+c y$ are all permutation polynomials modulo $n$.

Remark 1.2 It must be noted that $P(x, y)=a+b x+c y$ represents a groupoid over $\mathbb{Z}_{n} . P(x, y)$ represents a quasigroup over $\mathbb{Z}_{n}$ if and only if $\left(\mathbb{Z}_{n}, P\right)$ is a T-quasigroup. Hence whenever $\left(\mathbb{Z}_{n}, P\right)$ is a groupoid and not a quasigroup, $\left(\mathbb{Z}_{n}, P\right)$ is neither a T-quasigroup nor left linear quasigroup nor right linear quasigroup. Thus, the present study considers both T-quasigroup and non-T-quasigroup.

Theorem 1.3(Vadiraja and Shankar [27]) If $P(x, y)$ is a bivariate polynomial having no cross term, then $P(x, y)$ gives a Latin square if and only if $P(x, 0)$ and $P(0, y)$ are permutation polynomials.

The authors were able to establish the fact that Rivest's result for a bivariate polynomial over $\mathbb{Z}_{n}$ when $n=2^{w}$ is true for a linear-bivariate polynomial over $\mathbb{Z}_{n}$ when $n \neq 2^{w}$. Although the result of Rivest was found not to be true for quadratic-bivariate polynomials over $\mathbb{Z}_{n}$ when $n \neq 2^{w}$ with the help of counter examples, nevertheless some of such squares can be forced to be Latin squares by deleting some equal numbers of rows and columns.

Furthermore, Vadiraja and Shankar [27] were able to find examples of pairs of orthogonal Latin squares generated by bivariate polynomials over $\mathbb{Z}_{n}$ when $n \neq 2^{w}$ which was found impossible by Rivest for bivariate polynomials over $\mathbb{Z}_{n}$ when $n=2^{w}$.

### 1.4 Some Important Results on Medial-Like Identities

Some important results which we would find useful in our study are stated below.
Theorem 1.4(Polonijo [21]) For any groupoid $(Q, \cdot)$, any two of the three identities (55), (65) and (66) imply the third one.

Theorem 1.5(Polonijo [21]) Let $(Q, \cdot)$ be a commutative groupoid. Then $(Q, \cdot)$ is palindromic. Furthermore, the constraints (55) and (65) are equivalent, i.e a commutative groupoid $(Q, \cdot)$ is internally medial if and only if it is externally medial.

Theorem 1.6(Polonijo [21]) For any quasigroup $(Q, \cdot)$ and $i \in\{1,2, \ldots, 6\}, C_{i}$ is valid if and only if the quasigroup is commutative.

Theorem 1.7(Polonijo [21]) For any quasigroup $(Q, \cdot)$ and $i \in\{1,2, \ldots, 14\}, C M_{i}$ holds if and only if the quasigroup is both commutative and internally medial.

Theorem 1.8(Polonijo [21]) For any quasigroup $(Q, \cdot)$ and $i \in\{1,2, \ldots, 14\}, C M_{i}$ is valid if and only if the quasigroup is both commutative and externally medial.

Theorem 1.9(Polonijo [21]) A quasigroup $(Q, \cdot)$ is palindromic if and only if there exists an automorphism $\alpha$ such that

$$
\alpha(x \cdot y)=y \cdot x \forall x, y \in Q
$$

holds.
It is important to study the characterization of varieties of groupoids and quasigroups represented by linear-bivariate polynomials over the ring $\mathbb{Z}_{n}$ even though very few of such have been sighted as examples in the past.

## $\S 2$ Main Results

Theorem 2.1 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\left\{\mathbb{Z}_{n}, \mathbb{Z}_{p}\right\}$ such that "HYPO" is true. $P(x, y)$ represents a "NAME" \{groupoid, quasigroup $\}\left\{\left(\mathbb{Z}_{n}, P\right),\left(\mathbb{Z}_{p}, P\right)\right\}$ over $\left\{\mathbb{Z}_{n}, \mathbb{Z}_{p}\right\}$ if and only if " $N$ and $S$ " is true. (Table 1)

Proof There are 66 identities for which the theorem above is true for in a groupoid or quasigroup. For the sake of space, we shall only demonstrate the proof for one identity for each category.
(A) Those identities which involve one element only on each side of the equality sign:

Lemma 2.1 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a unipotent groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $(b+c)(x-y)=0$ for all $x, y \in \mathbb{Z}_{n}$.

Proof $P(x, y)$ satisfies the unipotent law $\Leftrightarrow P(x, x)=P(y, y) \Leftrightarrow a+b x+c x=a+b y+c y$ $\Leftrightarrow a+b x-c x-a-b y-c y=0 \Leftrightarrow(b+c)(x-y)=0$ as required.

Lemma 2.2 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a unipotent quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $(b+c)(x-y)=0$ and $(b, n)=$ $(c, n)=1$ for all $x, y \in \mathbb{Z}_{n}$.

Proof This is proved by using Lemma 2.1 and Theorem 1.2.

Theorem 2.2 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a unipotent groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $b+c \equiv 0(\bmod n)$.

Proof This is proved by using Lemma 2.1.

Theorem 2.3 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n}$. $P(x, y)$ represents a unipotent quasigroupp $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $b+c \equiv 0(\bmod n)$ and $(b, n)=$ $(c, n)=1$.

Proof This is proved by using Lemma 2.2.
Example 2.1 $P(x, y)=5 x+y$ is a linear bivariate polynomial over $\mathbb{Z}_{6}$. $\left(\mathbb{Z}_{6}, P\right)$ is a unipotent groupoid over $\mathbb{Z}_{6}$.

Example 2.2 $P(x, y)=1+5 x+y$ is a linear bivariate polynomial over $\mathbb{Z}_{6} .\left(\mathbb{Z}_{6}, P\right)$ is a unipotent quasigroup over $\mathbb{Z}_{6}$.
(B) Those identities which involve two elements on one or both sides of the equality sign:

Lemma 2.3 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . \quad P(x, y)$ represents a Stein third groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $a(1+b+c)+x\left(b^{2}+c^{2}\right)+$ $y(2 b c-1)=0$ for all $x, y \in \mathbb{Z}_{n}$.

Proof $P(x, y)$ satisfies the Stein third law $\Leftrightarrow P[P(x, y), P(y, x)]=y \Leftrightarrow a(1+b+c)+x\left(b^{2}+\right.$ $\left.c^{2}\right)+y(2 b c-1)=0$ as required.

Lemma 2.4 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . \quad P(x, y)$ represents a Stein third quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $a(1+b+c)+x\left(b^{2}+c^{2}\right)+$ $y(2 b c-1)=0$ and $(b, n)=(c, n)=1$ for all $x, y \in \mathbb{Z}_{n}$.

Proof This is proved by using Lemma 2.3 and Theorem 1.2.
Theorem 2.4 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a Stein third groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $b^{2}+c^{2} \equiv 0(\bmod n), 2 b c \equiv$ $1(\bmod n)$ and $a=0$.

Proof This is proved by using Lemma 2.3.

Theorem 2.5 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a Stein third quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $b^{2}+c^{2} \equiv 0(\bmod n), 2 b c \equiv$ $1(\bmod n)$ and $a=0$.

Proof This is proved by using Lemma 2.4.

Theorem 2.6 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{p}$ such that $a \neq 0$. $P(x, y)$ represents a Stein third groupoid $\left(\mathbb{Z}_{p}, P\right)$ over $\mathbb{Z}_{p}$ if and only if $b^{2}+c^{2} \equiv 0(\bmod p)$ and $2 b c \equiv 1(\bmod p)$.

Proof This is proved by using Lemma 2.3.

Theorem 2.7 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{p}$ such that $a \neq 0$. $P(x, y)$ represents a Stein third quasigroup $\left(\mathbb{Z}_{p}, P\right)$ over $\mathbb{Z}_{p}$ if and only if $b^{2}+c^{2} \equiv 0(\bmod p)$ and $2 b c \equiv 1(\bmod p)$.

Proof This is proved by using Lemma 2.4.
Example 2.3 $P(x, y)=2 x+3 y$ is a linear bivariate polynomial over $\mathbb{Z}_{5} .\left(\mathbb{Z}_{5}, P\right)$ is a Stein third groupoid over $\mathbb{Z}_{5}$.

Example 2.4 $P(x, y)=2 x+3 y$ is a linear bivariate polynomial over $\mathbb{Z}_{5} .\left(\mathbb{Z}_{5}, P\right)$ is a Stein third quasigroup over $\mathbb{Z}_{5}$.
(C) Those identities which involve three distinct elements on one or both sides of the equality sign:

Lemma 2.5 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . \quad P(x, y)$ represents an Abel-Grassman groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $(x-z)\left(b-c^{2}\right)=0$ for all $x, z \in \mathbb{Z}_{n}$.

Proof $P(x, y)$ satisfies the Abel-Grassman law $\Leftrightarrow P[x, P(y, z)]=P[z, P(y, x)] \Leftrightarrow P(x, a+$ $b y+c z)=P(z, a+b y+c x) \Leftrightarrow a+b x+c(a+b y+c z)=a+b z+c(a+b y+c x) \Leftrightarrow(x-z)\left(b-c^{2}\right)=0$ as required.

Lemma 2.6 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . \quad P(x, y)$ represents an Abel-Grassman quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $(x-z)\left(b-c^{2}\right)=0$ and $(b, n)=(c, n)=1$. for all $x, y, z \in \mathbb{Z}_{n}$.

Proof This is proved by using Lemma 2.5 and Theorem 1.2.
Theorem 2.8 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents an Abel-Grassman groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $c^{2} \equiv b(\bmod n)$.

Proof This is proved by using Lemma 2.5.
Theorem 2.9 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents an Abel-Grassman quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $c^{2} \equiv b(\bmod n)$ and $(b, n)=(c, n)=1$.

Proof This is proved by using Lemma 2.6.
Example 2.5 $P(x, y)=2+4 x+2 y$ is a linear bivariate polynomial over $\mathbb{Z}_{6}$. $\left(\mathbb{Z}_{6}, P\right)$ is an Abel-Grassman groupoid over $\mathbb{Z}_{6}$.

Example 2.6 $P(x, y)=2+4 x+2 y$ is a linear bivariate polynomial over $\mathbb{Z}_{5} .\left(\mathbb{Z}_{5}, P\right)$ is an Abel-Grassman quasigroup over $\mathbb{Z}_{5}$.
(D) Those involving four elements:

Lemma 2.7 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . \quad P(x, y)$ represents an external medial groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $w\left(b^{2}-c^{2}\right)+z\left(c^{2}-b^{2}\right)=0$ for all $w, z \in \mathbb{Z}_{n}$.

Proof $P(x, y)$ satisfies the external medial law $\Leftrightarrow P[P(w, x), P(y, z)]=P[P(z, x), P(y, w)]$
$\Leftrightarrow a+b(a+b w+c x)+c(a+b y+c z)=a+b(a+b z+c x)+c(a+b y+c w) \Leftrightarrow w\left(b^{2}-c^{2}\right)+z\left(c^{2}-b^{2}\right)=0$ as required.

Lemma 2.8 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents an external medial quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $w\left(b^{2}-c^{2}\right)+z\left(c^{2}-b^{2}\right)=0$ and $(b, n)=(c, n)=1$ for all $w, z \in \mathbb{Z}_{n}$.

Proof This is proved by using Lemma 2.7 and Theorem 1.2.
Theorem 2.10 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} .\left(\mathbb{Z}_{n}, P\right)$ represents an external medial groupoid over $\mathbb{Z}_{n}$ if and only if $b^{2} \equiv c^{2}(\bmod n)$.

Proof This is proved by using Lemma 2.7.

Theorem 2.11 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} .\left(\mathbb{Z}_{n}, P\right)$ represents an external medial quasigroup over $\mathbb{Z}_{n}$ if and only if $b^{2} \equiv c^{2}(\bmod n)$ and $(b, n)=$ $(c, n)=1$.

Proof This is proved by using Lemma 2.8 and Theorem 1.2.
Example 2.7 $P(x, y)=4+2 x+2 y$ is a linear bivariate polynomial over $\mathbb{Z}_{6}$. $\left(\mathbb{Z}_{6}, P\right)$ is an external medial groupoid over $\mathbb{Z}_{6}$.

Example 2.8 $P(x, y)=2+8 x+y$ is a linear bivariate polynomial over $\mathbb{Z}_{9}$. $\left(\mathbb{Z}_{9}, P\right)$ is an external medial quasigroup over $\mathbb{Z}_{9}$.
(E) Those involving left or right inverse elements:

Lemma 2.9 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a cross inverse property groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $a(b c-1)+x\left(b^{2} c+\right.$ $1-b-b c)+c y(b c-1)=0$ for all $x, y \in \mathbb{Z}_{n}$.

Proof $P(x, y)$ satisfies the cross inverse property $\left.\Leftrightarrow P\left[P(x, y), x^{\rho}\right)\right]=y \Leftrightarrow P(a+b x+$ $\left.c y, x^{\rho}\right)=y \Leftrightarrow a+b(a+b x+c y)+c x^{\rho}=y \Leftrightarrow a(b c-1)+x\left(b^{2} c+1-b-b c\right)+c y(b c-1)=0$ as required.

Lemma 2.10 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a cross inverse property quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $a(b c-1)+x\left(b^{2} c+\right.$ $1-b-b c)+c y(b c-1)=0$ and $(b, n)=(c, n)=1$ for all $x, y, z \in \mathbb{Z}_{n}$.

Proof This is proved by using Lemma 2.9 and Theorem 1.2.

Theorem 2.12 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{p}$ such that $a \neq 0$. $P(x, y)$ represents a CIP quasigroup $\left(\mathbb{Z}_{p}, P\right)$ over $\mathbb{Z}_{p}$ if and only if bc $\equiv 1(\bmod p)$.

Proof This is proved by using Lemma 2.10.
Theorem 2.13 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n}$ such that
$a \neq 0$ and $c$ is invertible in $\mathbb{Z}_{n} . P(x, y)$ represents a CIP groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $b c \equiv 1(\bmod n)$.

Proof This is proved by using Lemma 2.9.

Theorem 2.14 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n}$ such that $a \neq 0, c$ is invertible in $\mathbb{Z}_{n}$ and $(b, n)=(c, n)=1 . P(x, y)$ represents a CIP quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if and only if $b c \equiv 1(\bmod n)$.

Proof This is proved by using Lemma 2.10.

Theorem 2.15 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n} . P(x, y)$ represents a CIP groupoid $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if $b c \equiv 1(\bmod n)$.

Proof This is proved by using Lemma 2.9.

Theorem 2.16 Let $P(x, y)=a+b x+c y$ be a linear bivariate polynomial over $\mathbb{Z}_{n}$ such that $(b, n)=(c, n)=1 . P(x, y)$ represents a CIP quasigroup $\left(\mathbb{Z}_{n}, P\right)$ over $\mathbb{Z}_{n}$ if $b c \equiv 1(\bmod n)$.

Proof This is proved by using Lemma 2.10.

Example 2.9 $P(x, y)=2+4 x+4 y$ is a linear bivariate polynomial over $\mathbb{Z}_{5} \cdot\left(\mathbb{Z}_{5}, P\right)$ is a cross inverse property groupoid over $\mathbb{Z}_{5}$.

Example $2.10 \quad P(x, y)=3+4 x+4 y$ is a linear bivariate polynomial over $\mathbb{Z}_{5}$. $\left(\mathbb{Z}_{5}, P\right)$ is a cross inverse property quasigroup over $\mathbb{Z}_{5}$.

Table 1. Varieties of groupoids and quasigroups generated by $P(x, y)$ over $\mathbb{Z}_{n}$

| S/N | NAME | G | Q | $\mathbb{Z}_{n}$ | $\mathbb{Z}_{p}$ | HYPO | N AND S | EXAMPLE |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | Idempotent | $\checkmark$ |  | $\checkmark$ |  |  | $b+c=1, a=0$ | $5 x+2 y, \mathbb{Z}_{6}$ |
| 2 | Unipotent | $\checkmark$ |  | $\checkmark$ |  |  | $b+c=0$ | $2+4 x+2 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b+c=0,(b, n)=(c, n)=1$ | $2+5 x+y, \mathbb{Z}_{6}$ |
| 3 | Commut | $\checkmark$ |  | $\checkmark$ |  |  | $b=c$ | $1+4 x+4 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b=c,(b, n)=(c, n)=1$ | $1+5 x+5 y, \mathbb{Z}_{6}$ |
| 4 | Sade Right | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=-1$ | $2+6 x+4 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=-1$ | $1+5 x+4 y, \mathbb{Z}_{7}$ |
| 5 | Sade Left | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $c=-1$ | $2+4 x+5 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $c=-1$ | $2+5 x+5 y, \mathbb{Z}_{7}$ |
| 6 | Right <br> Alternative | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $3+x+y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $3+x+y, \mathbb{Z}_{7}$ |
| 7 | Left <br> Alternative | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $2+x+y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $2+x+y, \mathbb{Z}_{7}$ |


| S/N | NAME | G | Q | $\mathbb{Z}_{n}$ | $\mathbb{Z}_{p}$ | HYPO | N AND S | EXAMPLE |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 8 | Medial <br> Alternative | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $b \neq c$ | $b+c=1$ | $2+4 x+2 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $b \neq c$ | $b+c=1$ | $2+4 x+2 y, \mathbb{Z}_{7}$ |
| 9 | Right <br> Semi <br> Symmetry | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c=-1$ | $2+4 x+4 y, \mathbb{Z}_{5}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0$ | $b c=1, c^{2}=-b$ | $5 x+2 y, \mathbb{Z}_{9}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c=-1$ | $2+4 x+4 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0$ | $b c=1, c^{2}=-b$ | $5 x+2 y, \mathbb{Z}_{9}$ |
| 10 | Left <br> Semi <br> Symmetry | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c=-1$ | $3+4 x+4 y, \mathbb{Z}_{5}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0$ | $b=1, b^{2}=-c$ | $x+9 y, \mathbb{Z}_{10}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c=-1$ | $3+4 x+4 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0$ | $b=1, b^{2}=-c$ | $x+9 y, \mathbb{Z}_{10}$ |
| 11 | Stein First | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c$ | $3+4 x+4 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{5}$ |
| 12 | Stein <br> Second | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c$ | $3+4 x+4 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{5}$ |
| 13 | Schroder <br> Second | $\checkmark$ |  | $\checkmark$ |  |  | $b^{2}+c^{2}=1,2 b c=a=0$ | $2 x+3 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b^{2}+c^{2}=(b, n)=(c, n)=1,2 b c=a=0$ | ? |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b+c=-1, b^{2}+c^{2}=1,2 b c=0$ | ? |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b+c=-1, b^{2}+c^{2}=1,2 b c=0$ | ? |
| 14 | Stein Third | $\checkmark$ |  | $\checkmark$ |  |  | $b^{2}+c^{2}=0,2 b c=1, a=0$ | ? |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $(b, n)=(c, n)=2 b c=1, b^{2}+c^{2}=a=0$ | $?$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b^{2}+c^{2}=0,2 b c=1$, | $3+2 x+4 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b^{2}+c^{2}=0,2 b c=1$, | $2+2 x+4 y, \mathbb{Z}_{5}$ |
| 15 | Associative | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $2+x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $2+x+y, \mathbb{Z}_{6}$ |
| 16 | Slim | $\checkmark$ |  | $\checkmark$ |  | $a=0, c$ invert | $b c=0, c=1$ | ! |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, c$ invert | $b c=0, c=1,(b, n)=(c, n)=1$ | ? |
| 17 | Cyclic Associativity | $\checkmark$ |  | $\checkmark$ |  |  | $b=c=1$ | $3+x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b=c=1,(b, n)=(c, n)=1$ | $3+x+y, \mathbb{Z}_{6}$ |
| 18 | Right <br> Permutability | $\checkmark$ |  | $\checkmark$ |  |  | $b=1$ | $1+x+5 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b=1,(b, n)=(c, n)=1$ | $1+x+5 y, \mathbb{Z}_{6}$ |
| 19 | LeftPermutability | $\checkmark$ |  | $\checkmark$ |  |  | $c=1$ | $1+5 x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $c=1,(b, n)=(c, n)=1$ | $3+5 x+y, \mathbb{Z}_{6}$ |
| 20 | Abel <br> Grassman | $\checkmark$ |  | $\checkmark$ |  |  | $c^{2}=b$ | $2+4 x+2 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $c^{2}=b,(b, n)=(c, n)=1$ | $2+4 x+2 y, \mathbb{Z}_{9}$ |
| 21 | Commuting Product | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $1+x+y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $1+x+y, \mathbb{Z}_{7}$ |
| 22 | Dual Comm Product | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $1+x+y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c=1$ | $1+x+y, \mathbb{Z}_{7}$ |
| 23 | Right Transitivity | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=1, c=-1$ | $2+x+6 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=1, c=-1$ | $2+x+6 y, \mathbb{Z}_{7}$ |
| 24 | Left <br> Transitivity | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=-1, c=1$ | $2+6 x+y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=-1, c=1$ | $2+6 x+y, \mathbb{Z}_{7}$ |
| 25 | Schweitzer | $\checkmark$ |  | $\checkmark$ |  | $b, c$ invert | $b=1, c=-1$ | $2+x+5 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $b, c$ invert | $b=1, c=-1,(b, n)=(c, n)=1$ | $2+x+5 y, \mathbb{Z}_{6}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=1, c=-1$ | $3+x+6 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=1, c=-1$ | $3+x+6 y, \mathbb{Z}_{7}$ |
| 26 | Dual of Schweitzer | $\checkmark$ |  | $\checkmark$ |  | $b, c$ invert | $b=1, c=-1$ | $2+x+5 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $b, c$ invert | $b=1, c=-1,(b, n)=(c, n)=1$ | $2+x+5 y, \mathbb{Z}_{6}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=1, c=-1$ | $3+x+6 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=1, c=-1$ | $3+x+6 y, \mathbb{Z}_{7}$ |


| S/N | NAME | G | Q | $\mathbb{Z}_{n}$ | $\mathbb{Z}_{p}$ | HYPO | N AND S | EXAMPLE |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 27 | Right Self Distributive | $\checkmark$ |  |  | $\checkmark$ |  | $c=1-b, a=0$ | $3 x+5 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  | $c=1-b, a=0$ | $3 x+5 y, \mathbb{Z}_{7}$ |
| 28 | Left Self Distributive | $\checkmark$ |  |  | $\checkmark$ |  | $c=1-b, a=0$ | $3 x+5 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  | $c=1-b, a=0$ | $3 x+5 y, \mathbb{Z}_{7}$ |
| 29 | Right <br> Abelian <br> Distributivity | $\checkmark$ |  | $\checkmark$ |  | $b, c$ invert | $b=c, 2 b^{2}=b$ | ? |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $b, c$ invert | $b=c, 2 b^{2}=b$ | ? |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a \neq 0$ | $b=c, 2 b^{2}=b$ | ? |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a \neq 0$ | $b=c, 2 b^{2}=b$ | ? |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c, 2 b=1$ | $2+3 x+3 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c, 2 b=1$ | $2+3 x+3 y, \mathbb{Z}_{5}$ |
| 30 | Left <br> Abelian <br> Distributivity | $\checkmark$ |  | $\checkmark$ |  | $b, c$ invert | $b=c, 2 b^{2}=b$ |  |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $b, c$ invert | $b=c, 2 b^{2}=b$ | ? |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a \neq 0$ | $b=c, 2 b^{2}=b$ | ? |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a \neq 0$ | $b=c, 2 b^{2}=b$ | ? |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b=c, 2 b=1$ | $2+3 x+3 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b=c, 2 b=1$ | $2+3 x+3 y, \mathbb{Z}_{5}$ |
| 31 | Bol <br> Moufang | $\checkmark$ |  | $\checkmark$ |  |  | $b=c=1$ | $2+x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $(b, n)=(c, n)=1$ | $b=c=1$ | $2+x+y, \mathbb{Z}_{6}$ |
| 32 | Dual Bol Moufang | $\checkmark$ |  | $\checkmark$ |  |  | $b=c=1$ | $2+x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $(b, n)=(c, n)=1$ | $b=c=1$ | $2+x+y, \mathbb{Z}_{6}$ |
| 33 | Moufang | $\checkmark$ |  |  | $\checkmark$ |  | $b=c=1, a=0$ | $x+y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  | $b=c=1, a=0$ | $x+y, \mathbb{Z}_{5}$ |
| 34 | R Bol | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b^{2}=1, b=c=1$ | $2+x+y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b^{2}=1, b=c=1$ | $2+x+y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $-1 \neq b \neq c$ | $b^{2}=1, c=1, a=0$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $-1 \neq b \neq c$ | $b^{2}=1, c=1, a=0$ | $8 x+y, \mathbb{Z}_{63}$ |
| 35 | L Bol | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $c^{2}=1, b=c=1$ | $2+x+y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $c^{2}=1, b=c=1$ | $2+x+y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $-1 \neq b \neq c$ | $c^{2}=1, b=1, a=0$ | $x+8 y, \mathbb{Z}_{63}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $-1 \neq b \neq c$ | $c^{2}=1, b=1, a=0$ | $x+8 y, \mathbb{Z}_{63}$ |
| 36 | $\mathrm{RC}_{4}$ | $\checkmark$ |  |  | $\checkmark$ | $a=0$ | $c=b^{2}=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a=0$ | $c=b^{2}=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, b, c$ invert | $c=b^{2}=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, b, c$ invert | $c=b^{2}=(b, n)=(c, n)=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  |  | $b=-1, c=1$ | $2+5 x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b=-1, c=(b, n)=(c, n)=1$ | $2+5 x+y, \mathbb{Z}_{6}$ |
| 37 | $\mathrm{LC}_{4}$ | $\checkmark$ |  |  | $\checkmark$ | $a=0$ | $b=c^{2}=1$ | $x+8 y, \mathbb{Z}_{63}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a=0$ | $b=c^{2}=1$ | $x+8 y, \mathbb{Z}_{63}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, b, c$ invert | $b=c^{2}=1$ | $x+3 y, \mathbb{Z}_{8}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, b, c$ invert | $b=c^{2}=(b, n)=(c, n)=1$ | $x+4 y, \mathbb{Z}_{15}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  |  | $b=-1, c=1$ | $2+5 x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b=-1, c=(b, n)=(c, n)=1$ | $2+5 x+y, \mathbb{Z}_{6}$ |
| 38 | $\mathrm{RC}_{1}$ | $\checkmark$ |  |  | $\checkmark$ | $a=0$ | $c=b^{2}=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a=0$ | $c=b^{2}=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, b, c$ invert | $c=b^{2}=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, b, c$ invert | $c=b^{2}=(b, n)=(c, n)=1$ | $8 x+y, \mathbb{Z}_{63}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  |  | $b=-1, c=1$ | $2+5 x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b=-1, c=(b, n)=(c, n)=1$ | $2+5 x+y, \mathbb{Z}_{6}$ |


| S/N | NAME | G | Q | $\mathbb{Z}_{n}$ | $\mathbb{Z}_{p}$ | HYPO | N AND S | EXAMPLE |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 39 | $\mathrm{LC}_{1}$ | $\checkmark$ |  |  | $\checkmark$ | $a=0, c \neq 1$ | $c=-1$ | $3 x+6 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a=0, c \neq 1$ | $c=-1$ | $3 x+6 y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, c \neq 1, c$ invert | $c=-1$ | $5 x+5 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, c \neq 1, c$ invert | $c=-1,(b, n)=(c, n)=1$ | $5 x+5 y, \mathbb{Z}_{6}$ |
| 40 | $\mathrm{LC}_{3}$ | $\checkmark$ |  | $\checkmark$ |  |  | $c=1, \quad b=-2$ | $3+4 x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $c=1, b=-2,(b, n)=(c, n)=1$ | $2+5 x+y, \mathbb{Z}_{7}$ |
| 41 | $\mathrm{RC}_{3}$ | $\checkmark$ |  | $\checkmark$ |  |  | $c=1, \quad b=-2$ | $3+4 x+y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $c=1, b=-2,(b, n)=(c, n)=1$ | $2+5 x+y, \mathbb{Z}_{7}$ |
| 42 | C-Law | $\checkmark$ |  |  | $\checkmark$ | $a=0$ | $b=c=-1$ | $4 x+4 y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a=0$ | $b=c=-1$ | $4 x+4 y, \mathbb{Z}_{5}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a \neq 0, b \neq 1, b, c$ inv | $b=c=-1$ | $3+5 x+5 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a \neq 0, b \neq 1, b, c$ inv | $b=c=-1,(b, n)=(c, n)=1$ | $3+5 x+5 y, \mathbb{Z}_{6}$ |
| 43 | LIP | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $c^{2}=b^{2}=b c=1$ | ? |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $c^{2}=b^{2}=b c=1$ | ? |
| 44 | RIP | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $c^{2}=b^{2}=b c=1$ | ? |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $c^{2}=b^{2}=b c=1$ | ? |
| 45 | 1st Right CIP | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b c=1$ | $2+3 x+4 y, \mathbb{Z}_{11}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b c=1$ | $2+3 x+4 y, \mathbb{Z}_{11}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a \neq 0, c$ inv | $b c=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a \neq 0, c$ inv | $b c=1,(b, n)=(c, n)=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
| 46 | 2nd Right CIP | $\checkmark$ |  | $\checkmark$ |  |  | $b c=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b c=1,(b, n)=(c, n)=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
| 47 | 1st Left CIP | $\checkmark$ |  |  | $\checkmark$ | $a \neq 0$ | $b c=1$ | $2+3 x+4 y, \mathbb{Z}_{11}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a \neq 0$ | $b c=1$ | $2+3 x+4 y, \mathbb{Z}_{11}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a \neq 0, b$ inv | $b c=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a \neq 0, b$ inv | $b c=1,(b, n)=(c, n)=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
| 48 | 2nd Left CIP | $\checkmark$ |  | $\checkmark$ |  |  | $b c=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  | $b c=1,(b, n)=(c, n)=1$ | $3+3 x+3 y, \mathbb{Z}_{8}$ |
| 49 | R AAIP | $\checkmark$ |  |  | $\checkmark$ | $b c+b \neq 1$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{11}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $b c+b \neq 1$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{11}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $c \neq b$ | $b+b c=1$ | $2+3 x+y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $c \neq b$ | $b+b c=1$ | $2+3 x+y, \mathbb{Z}_{5}$ |
| 50 | L AAIP | $\checkmark$ |  |  | $\checkmark$ | $b c+b \neq 1$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{11}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $b c+b \neq 1$ | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{11}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ | $c \neq b$ | $b+b c=1$ | $2+3 x+y, \mathbb{Z}_{5}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $c \neq b$ | $b+b c=1$ | $2+3 x+y, \mathbb{Z}_{5}$ |
| 51 | R AIP | $\checkmark$ |  | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 52 | L AIP | $\checkmark$ |  | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 53 | R SAIP | $\checkmark$ |  | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |


| S/N | NAME | G | Q | $\mathbb{Z}_{n}$ | $\mathbb{Z}_{p}$ | HYPO | N AND S | EXAMPLE |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 54 | L SAIP | $\checkmark$ |  | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 55 | R WIP | $\checkmark$ |  |  | $\checkmark$ | $a=0, c^{2} \neq 0$ | $b c=1$ | $3 x+5 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a=0, c^{2} \neq 0$ | $b c=1$ | $3 x+5 y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, c$ inv | $b c=1$ | $3 x+4 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, c$ inv | $b c=1,(b, n)=(c, n)=1$ | ? |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, b c+b \neq 1$ | $b c=1$ | ? |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, b c+b \neq 1$ | $b c=1,(b, n)=(c, n)=1$ | ? |
| 56 | L WIP | $\checkmark$ |  |  | $\checkmark$ | $a=0, b^{2} \neq 0$ | $b c=1$ | $3 x+5 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ | $a=0, b^{2} \neq 0$ | $b c=1$ | $3 x+5 y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, b$ inv | $b c=1$ | $3 x+4 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, b$ inv | $b c=1,(b, n)=(c, n)=1$ | ? |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $a=0, b c+c \neq 1$ | $b c=1$ | ? |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $a=0, b c+c \neq 1$ | $b c=1,(b, n)=(c, n)=1$ | ? |
| 57 | $\mathrm{E}_{1}$ | $\checkmark$ | $\checkmark$ |  |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 58 | $\mathrm{E}_{\mathrm{r}}$ | $\checkmark$ | $\checkmark$ |  |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 59 | Right F | $\checkmark$ | $\checkmark$ |  |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 60 | Left F | $\checkmark$ | $\checkmark$ |  |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 61 | Medial | $\checkmark$ | $\checkmark$ |  |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 62 | Specialized <br> Medial | $\checkmark$ | $\checkmark$ |  |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  | $\checkmark$ |  |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  |  | $a+b x+c y, \mathbb{Z}_{n}$ |
| 63 | First <br> Rectangle | $\checkmark$ |  |  | $\checkmark$ |  | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $c$ inv | $b=c$ | $2+4 x+4 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $c$ inv | $b=c,(b, n)=(c, n)=1$ | $2+4 x+4 y, \mathbb{Z}_{6}$ |
| 64 | Second <br> Rectangle | $\checkmark$ |  |  | $\checkmark$ |  | $b=-c$ | $2+4 x+4 y, \mathbb{Z}_{7}$ |
|  |  |  | $\checkmark$ |  | $\checkmark$ |  | $b=-c$ | $2+4 x+4 y, \mathbb{Z}_{7}$ |
|  |  | $\checkmark$ |  | $\checkmark$ |  | $b$ inv | $b=-c$ | $2+4 x+4 y, \mathbb{Z}_{6}$ |
|  |  |  | $\checkmark$ | $\checkmark$ |  | $b$ inv | $b=-c,(b, n)=(c, n)=1$ | $2+4 x+4 y, \mathbb{Z}_{6}$ |
| 65 | $\mathrm{C}_{i}, i=1-6$ |  | $\checkmark$ |  | $\checkmark$ |  | $b=c$ | $3+5 x+5 y, \mathbb{Z}_{7}$ |
| 66 | $\mathrm{CM}_{i}, i=1-14$ |  | $\checkmark$ |  | $\checkmark$ | $b \neq-c$ | $b=c$ | $3+5 x+5 y, \mathbb{Z}_{7}$ |

Remark 2.1 A summary of the results on the characterization of groupoids and quasigroups generated by $P(x, y)$ is exhibited in Table 1. In this table, $G$ stands for groupoid, $Q$ stands for quasigroup, HYPO stands for hypothesis, N AND S stands for necessary and sufficient condition(s). Cells with question marks mean examples could not be gotten.

## References

[1] R. H.Bruck (1966), A Survey of Binary Systems, Springer -Verlag, Berlin- GöttingenHeidelberg, 185pp.
[2] O.Chein, H. O.Pflugfelder and J. D. H. Smith (1990), Quasigroups and Loops: Theory and Applications, Heldermann Verlag, 568pp.
[3] J.Dene and A.D.Keedwell (1974), Latin Squares and Their Applications, the English University press Lts, 549pp.
[4] E.G.Goodaire, E.Jespers and C.P.Milies (1996), Alternative Loop Rings, NHMS(184), Elsevier, 387 pp .
[5] T.G.Jaiyeola (2009), A Study of New Concepts in Smarandache Quasigroups and Loops, ProQuest Information and Learning(ILQ), Ann Arbor, USA, 127pp.
[6] A.D.Keedwell (2008), When is it hard to show that a quasigroup is a loop?, Comment. Math. Carolin., Vol. 49, No. 2, 241-247.
[7] A.D.Keedwell (2009), Realizations of loops and groups defined by short identities, Comment. Math. Carolin., Vol. 50, No. 3, 373-383.
[8] A.D.Keedwell and V.A.Shcherbacov (2003), Construction and properties of $(r, s, t)$-inverse quasigroups (I), Discrete Math., 266, 275-291.
[9] A.D.Keedwell and V.A.Shcherbacov (2004), Construction and properties of $(r, s, t)$-inverse quasigroups (II), Discrete Math., 288, 61-71.
[10] M.K.Kinyon and J. D.Phillips (2004), Axioms for trimedial quasigroups, Comment. Math. Univ. Carolinae, 45, 287-294.
[11] M.K.Kinyon and J.D.Phillips (2002), A note on trimedial quasigroups, Quasigroups and Related Systems, Vol.9, 65-66.
[12] M.K.Kinyon and J.D.Phillips (2005), Rectangular loops and rectangular quasigroups, Comput. Math. Appl., 49, 1679-1685.
[13] O.U.Kirnasovsky (1995), Linear isotopes of small order groups, Quasigroups and Related Systems, Vol. 2, 51-82.
[14] V.Krcadinac, V.Volenec (2005), A class of quasigroups associated with a cubic Pisot number, Quasigroups and Related Systems, 13 (2005), No. 2, 269-280.
[15] K.Kunen (1996), Quasigroups, loops and associative laws, J. Alg., 185, 194-204.
[16] K.Kunen (1996), Moufang Quasigroups, J. Alg., 183, 231-234.
[17] R.A.Mollin, C. Small (1987), On permutation polynomials over finite fields, Internat. J. Math. and Math. Sci., Vol. 10, No 3, 535-544.
[18] D.C.Murdoch (2001), Quasigroups which satisfy certain generalized associative laws, Amer. J. Math., 61 (1939), No. 2, 509-522.
[19] H.O.Pflugfelder (1990), Quasigroups and Loops: Introduction, Sigma Series in Pure Math. 7, Heldermann Verlag, Berlin, 147pp.
[20] M.J.Pelling, D.G.Rogers (1979), Stein quasigroups (II): Algebraic aspects, Bull. Austral. Math. Soc., 20, 321-344.
[21] M.Polonijo (2005), On medial-like identities, Quasigroups and Related Systems, 13, 281288.
[22] R.L.Rivest (2001), Permutation polynomials Modulo $2^{w}$, Finite Fields and Their Applications, 7, 287-292.
[23] L.Rudolf, G.L.Mullen (1988), When does a polynomial over a finite field Permute the elements of the field?, The American Mathematical Monthly, Vol. 95, No. 3, 243-246.
[24] L.Rudolf, G.L.Mullen (1993), When does a polynomial over a finite field Permute the elements of the field? (II), The American Mathematical Monthly, Vol. 100, No. 1, 71-74.
[25] L.V.Sabinin (1999), Smooth Quasigroups and Loops, Kluver Academic Publishers, Dordrecht, 249pp.
[26] J.D.H.Smith (2007), An Introduction to Quasigroups and Their Representations, Taylor and Francis Group, LLC.
[27] G.R.Vadiraja Bhatta and B.R.Shankar (2009), Permutation polynomials modulo $n, n \neq 2^{w}$ and Latin squares, International J. Math. Combin., 2, 58-65.
[28] W.B.Vasantha Kandasamy (2002), Smarandache Loops, Department of Mathematics, Indian Institute of Technology, Madras, India, 128pp.

# New Characterizations for Bertrand Curves 

 in Minkowski 3-SpaceBahaddin Bukcu<br>Gaziosmanpasa University, Faculty of Sciences and Arts<br>Department of Mathematics, Tasliciftlik Campus, 60250, Tokat-Turkey)<br>Murat Kemal Karacan<br>Usak University, Faculty of Sciences and Arts<br>Department of Mathematics, 1 Eylul Campus, 64200, Usak-Turkey<br>Nural Yuksel<br>Erciyes University, Faculty of Art and Sciences, Department of Mathematics, Kayseri -Turkey<br>E-mail: bbukcu@yahoo.com, murat.karacan@usak.edu.tr, yukseln@erciyes.edu.tr


#### Abstract

Bertrand curves have been investigated in Lorentzian and Minkowski spaces and some characterizations have been given in $[1,2,6]$. In this paper, we have investigated the relations between Frenet vector fields and curvatures and torsions of Bertrand curves at the corresponding points in Minkowski 3-space.


Key Words: Bertrand curves, constant curvature and torsion, Minkowski 3- Space.
AMS(2010): 53A04,53A35,53B30

## §1. Introduction

In the study of the fundamental theory and the characterizations of space curves, the corresponding relations between the curves are the very interesting and important problem. The well-known Bertrand curve is characterized as a kind of such corresponding relation between the two curves. J. Bertrand studied curves in Euclidean 3-space whose principal normals are the principal normals of another curve. Such a curve is nowadays called a Bertrand curve. Bertrand curves have a characteristic property that curvature and torsion are in linear relation.In the recent work [2], the authors studied spacelike and timelike Bertrand curves in Minkowski 3-space. (See also independently obtained results by [6]).

In this paper, we have investigated the relations between Frenet vector fields and curvatures and torsions of Bertrand curves at the corresponding points in Minkowski 3-space.

[^8]
## §2. Preliminaries

The Minkowski 3 -space $E_{1}^{3}$ is the Euclidean 3 -space $E^{3}$ provided with the standard flat metric given by

$$
\langle,\rangle=-d x_{1}+d x_{3}+d x_{3}
$$

where $\left(x_{1}, x_{2}, x_{3}\right)$ is a rectangular coordinate system of $E_{1}^{3}$. Since $\langle$,$\rangle is an indefinite metric,$ recall that a vector $v \in E_{1}^{3}$ can have one of three Lorentzian causal characters: it can be spacelike if $\langle v, v\rangle>0$ or $v=0$, timelike if $\langle v, v\rangle<0$ and null (lightlike) if $\langle v, v\rangle=0$ and $v \neq 0$. Similarly, an arbitrary curve $\alpha=\alpha(s)$ in $E_{1}^{3}$ can locally be spacelike, timelike or null (lightlike), if all of its velocity vectors $\alpha^{\prime}(s)$ are respectively spacelike, timelike or null (lightlike).

Minkowski space is originally from the relativity in Physics. In fact, a timelike curve corresponds to the path of an observer moving at less than the speed of light. Denote by $\{T, N, B\}$ the moving Frenet frame along the curve $\alpha(s)$ in the space $E_{1}^{3}$. For an arbitrary curve $\alpha(s)$ in the space $E_{1}^{3}$, the following Frenet formulae are given. If $\alpha$ is timelike curve, then the Frenet formulae read

$$
\left[\begin{array}{c}
T^{\prime}  \tag{1.1}\\
N^{\prime} \\
B^{\prime}
\end{array}\right]=\left[\begin{array}{ccc}
0 & \kappa & 0 \\
\kappa & 0 & \tau \\
0 & -\tau & 0
\end{array}\right]\left[\begin{array}{l}
T \\
N \\
B
\end{array}\right]
$$

where $\langle T, T\rangle_{L}=-1,\langle N, N\rangle_{L}=1,\langle B, B\rangle_{L}=1,\langle T, N\rangle_{L}=\langle N, B\rangle_{L}=\langle T, B\rangle_{L}=0$. If $\alpha$ is a spacelike curve with a spacelike principal normal, then the Frenet formulae read

$$
\left[\begin{array}{c}
T^{\prime}  \tag{1.2}\\
N^{\prime} \\
B^{\prime}
\end{array}\right]=\left[\begin{array}{ccc}
0 & \kappa & 0 \\
-\kappa & 0 & \tau \\
0 & \tau & 0
\end{array}\right]\left[\begin{array}{l}
T \\
N \\
B
\end{array}\right]
$$

where $\langle T, T\rangle_{L}=\langle N, N\rangle_{L}=1,\langle B, B\rangle_{L}=-1,\langle T, N\rangle_{L}=\langle N, B\rangle_{L}=\langle T, B\rangle_{L}=0$.If $\alpha$ is a spacelike curve with a spacelike binormal, then the Frenet formulae read

$$
\left[\begin{array}{c}
T^{\prime}  \tag{1.3}\\
N^{\prime} \\
B^{\prime}
\end{array}\right]=\left[\begin{array}{lll}
0 & \kappa & 0 \\
\kappa & 0 & \tau \\
0 & \tau & 0
\end{array}\right]\left[\begin{array}{l}
T \\
N \\
B
\end{array}\right]
$$

where $\langle T, T\rangle_{L}=\langle B, B\rangle_{L}=1,\langle N, N\rangle_{L}=-1,\langle T, N\rangle_{L}=\langle N, B\rangle_{L}=\langle T, B\rangle_{L}=0[4,7,11]$.

## §3. Bertrand Curves in Minkowski 3-Space

Definition 3.1 $([1,2,6])$ Let $\beta_{1}$ and $\beta_{2}$ be two unit speed regular curves in $E_{1}^{3}$, and $\left\{T_{1}, N_{1}, B_{1}\right\}$ and $\left\{T_{2}, N_{2}, B_{2}\right\}$ also be Frenet Frames on these curves, respectively. $\beta_{1}$ and $\beta_{2}$ are called Bertrand curves if $N_{1}$ and $N_{2}$ are linearly dependent. We say that $\beta_{2}$ is a Bertrand mate for $\beta_{1}$ and $\beta_{2}$ are Bertrand curves. And $\left(\beta_{1}, \beta_{2}\right)$ is called a Bertrand couple and we can write

$$
\beta_{2}(s)=\beta_{1}(s)+r N_{1}(s) .
$$

Theorem 3.1 If there exists a one-to-one correspondence between the points of the spacelike curves $C_{1}$ and $C_{2}$ with timelike principal normal, such that at corresponding points $P_{1}$ on $C_{1}$ and $P_{2}$ on $C_{2}$, then the following statements hold:
(1) The curvature $\kappa_{1}$ of $C_{1}$ is a constant;
(2) The torsion $\tau_{2}$ of $C_{2}$ is constant;
(3) The unit tangent vector $T_{1}$ of $C_{1}$ is parallel to the unit tangent vector $T_{2}$ of $C_{2}$.

Then the curve $C$ generated by $P$ that divides the segment $P_{1} P_{2}$ in ratio $m: 1$ is a spacelike Bertrand curve with timelike principal normal.

Proof We shall use the subscripts 1, 2 to designate the geometric quantites corresponding to the curves $C_{1}, C_{2}$ while the same letters without subscripts will refer to the spacelike curve $C$ with timelike principal normal.

Let $\alpha(s), \alpha_{1}(s), \alpha_{2}(s)$ be the coordinat vectors at the points $P, P_{1}, P_{2}$ on the curves $C, C_{1}, C_{2}$ respectively. Then the from convex combination of points $P_{1}$ and $P_{2}$, the equation of point $P$ is

$$
\begin{equation*}
\alpha(s)=m \alpha_{1}(s)+(1-m) \alpha_{2}(s), m \in R \tag{2.1}
\end{equation*}
$$

while by hypothesis,

$$
\begin{equation*}
\left\|T_{1}\right\|=\left\|T_{2}\right\|=1, \quad T_{1}=T_{2} \tag{2.2}
\end{equation*}
$$

On differentiating Eq.(2.1) we have

$$
\begin{equation*}
T d s=m T_{1} d s_{1}+(1-m) T_{2} d s_{2}=\left(m d s_{1}+(1-m) d s_{2}\right) T_{1} \tag{2.3}
\end{equation*}
$$

which shows that $T$ is parallel to $T_{1}$ and $T_{2}$ and always can be chosen so that

$$
\begin{equation*}
T=T_{1}=T_{2} \tag{2.4}
\end{equation*}
$$

and

$$
\begin{equation*}
d s=m d s_{1}+(1-m) d s_{2} \tag{2.5}
\end{equation*}
$$

Differentiating of Eq.(2.4) gives

$$
\begin{equation*}
\kappa N d s=\kappa_{1} N_{1} d s_{1}=\kappa_{2} N_{2} d s_{2} \tag{2.6}
\end{equation*}
$$

and if we assume that $\kappa, \kappa_{1}, \kappa_{2}$ are positive, then

$$
\begin{equation*}
N=N_{1}=N_{2} \tag{2.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\kappa d s=\kappa_{1} d s_{1}=\kappa_{2} d s_{2} \tag{2.8}
\end{equation*}
$$

From Eq.(2.4) and Eq.(2.7)

$$
\begin{equation*}
B=B_{1}=B_{2} \tag{2.9}
\end{equation*}
$$

and differentiating

$$
\begin{equation*}
\tau N d s=\tau_{1} N_{1} d s_{1}=\tau_{2} N_{2} d s_{2} \tag{2.10}
\end{equation*}
$$

Elimination of $d s, d s_{1}, d s_{2}$ gives

$$
\left(\frac{m}{\kappa_{1}}\right) \kappa+\left(\frac{1-m}{\tau_{2}}\right) \tau=1 ; \kappa_{1} \neq 0, \tau_{2} \neq 0
$$

which is the desired result, since $m, \kappa_{1}, \tau_{2}$ are constant. If instead of $T_{1}=T_{2}$ were given the condition $B_{1}=B_{2}$, the same result would follow in the same manner.

Theorem 3.2 If condition (c) of Theorem 3.1 is modifed so that at corresponding points $P_{1}$ and $P_{2}$, the binormals $B_{1}$ and $B_{2}$ are parallel, then the curve $C$ is a spacelike Bertrand curve with timelike principal normal.

Proof Since $B_{1}=B_{2}$ then

$$
\begin{equation*}
\tau_{1} N_{1}=\tau_{2} N_{2} \frac{d s_{2}}{d s_{1}} \tag{2.11}
\end{equation*}
$$

where $N_{1}$ and $N_{2}$ are the unit normal vectors of $\alpha_{1}$ and $\alpha_{2}$ at the points $P_{1}$ and $P_{2}$ with arclength parametrization. Hence $N_{1}=N_{2}$ and $T_{1}=N_{1} \times B_{1}=N_{2} \times B_{2}$, we know $T_{1}$ parallel to $T_{2}$. By Theorem 2.1, $C$ is a spacelike Bertrand curve with timelike principal normal.

Theorem 3.3 If condition (c) of Theorem 3.1 is modified so that at corresponding points $P_{1}$ and $P_{2}$ the tangent at $P_{1}$ is parallel to the binormal $B_{2}$ at $P_{2}$, then the curve $C$ is a spacelike Bertrand curve with timelike principal normal.

Proof Since $T_{1}=B_{2}$, it follow that

$$
\begin{equation*}
\kappa_{1} N_{1}=\tau_{2} N_{2} \frac{d s_{2}}{d s_{1}} \tag{2.12}
\end{equation*}
$$

Hence $N_{1}$ is parallel to $N_{2}$ and since $N_{1}$ and $N_{2}$ are unit vectors,

$$
\begin{equation*}
N_{1}=N_{2} \tag{2.13}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d s_{2}}{d s_{1}}=\frac{\kappa_{1}}{\tau_{2}} \tag{2.14}
\end{equation*}
$$

since $B_{1}=T_{1} \times N_{1}=B_{2} \times N_{2}=-T_{2}$, we have

$$
\begin{equation*}
B_{1}=-T_{2} \tag{2.15}
\end{equation*}
$$

Let $\alpha, \alpha_{1}, \alpha_{2}$ be the coordinate vectors at the points $P, P_{1}, P_{2}$ on the curves $C, C_{1}, C_{2}$, respectively. Then

$$
\begin{align*}
\alpha= & m \alpha_{1}+(1-m) \alpha_{2}  \tag{2.16}\\
\frac{d \alpha}{d s} & =m \frac{d \alpha_{1}}{d s}+(1-m) \frac{d \alpha_{2}}{d s} \\
\frac{d \alpha}{d s} & =m \frac{d \alpha_{1}}{d s_{1}} \frac{d s_{1}}{d s}+(1-m) \frac{d \alpha_{2}}{d s_{2}} \frac{d s_{2}}{d s} \\
T & =\left(m T_{1}+(1-m) \frac{d s_{2}}{d s_{1}} T_{2}\right) \frac{d s_{1}}{d s} \\
& =\left(m T_{1}+\frac{\kappa_{1}}{\tau_{2}}(1-m) T_{2}\right) \frac{d s_{1}}{d s}
\end{align*}
$$

$$
\begin{equation*}
T=m_{1} T_{1}+m_{2} T_{2}, \tag{2.17}
\end{equation*}
$$

where

$$
\begin{aligned}
& m_{1}=m \frac{d s_{1}}{d s}=\frac{m \tau_{2}}{\sqrt{\left(m \tau_{2}\right)^{2}+\left[\kappa_{1}(1-m)\right]^{2}}}, m_{1}=\text { const } . \\
& m_{2}=(1-m) \frac{d s_{2}}{d s}=\frac{(1-m) \kappa_{1}}{\sqrt{\left(m \tau_{2}\right)^{2}+\left[\kappa_{1}(1-m)\right]^{2}}}, m_{2}=\text { const. }
\end{aligned}
$$

Differentiating Eq.(2.17), one gets

$$
\begin{equation*}
\kappa N=\kappa_{1} m_{1} N_{1} \frac{d s_{1}}{d s}=\kappa_{2} m_{2} N_{2} \frac{d s_{2}}{d s} \tag{2.18}
\end{equation*}
$$

Hence

$$
\begin{equation*}
N=N_{1}=N_{2} \tag{2.19}
\end{equation*}
$$

and

$$
\begin{equation*}
\kappa=\kappa_{1} m_{1} \frac{d s_{1}}{d s}+\kappa_{2} m_{2} \frac{d s_{2}}{d s} \tag{2.20}
\end{equation*}
$$

Using Eq.(2.7) and Eq.(2.9), one finds that

$$
\begin{equation*}
B=m_{1} B_{1}+m_{2} B_{2} \tag{2.21}
\end{equation*}
$$

Differentiating Eq.(2.11), one gets

$$
\begin{equation*}
\tau N=\tau_{1} m_{1} \frac{d s_{1}}{d s} N_{1}+\tau_{2} m_{2} \frac{d s_{2}}{d s} N_{2} \tag{2.22}
\end{equation*}
$$

Hence

$$
\begin{equation*}
\tau=\tau_{1} m_{1} \frac{d s_{1}}{d s}+\tau_{2} m_{2} \frac{d s_{2}}{d s} \tag{2.23}
\end{equation*}
$$

Using Eq.(2.14) and Eq.(2.15), one gets

$$
\begin{equation*}
\frac{d s_{2}}{d s_{1}}=-\frac{\tau_{1}}{\kappa_{2}}=\frac{\kappa_{1}}{\tau_{2}} \tag{2.24}
\end{equation*}
$$

and

$$
\frac{\tau_{1}}{\kappa_{1}}=-\frac{\kappa_{2}}{\tau_{2}}
$$

Let

$$
M_{1}=m_{1} \frac{d s_{1}}{d s}, \quad M_{2}=m_{2} \frac{d s_{2}}{d s}
$$

Then using Eq.(2.20) and Eq.(2.23), one gets

$$
\begin{aligned}
\frac{\kappa}{M_{2} \tau_{2}}+\frac{\tau}{M_{1} \kappa_{1}} & =\frac{\kappa_{1}}{\tau_{2}} \frac{M_{1}}{M_{2}}+\left(\frac{\kappa_{2} M_{2}}{\tau_{2} M_{2}}+\frac{\tau_{1} M_{1}}{\kappa_{1} M_{1}}\right)+\frac{\tau_{2}}{\kappa_{1}} \frac{M_{2}}{M_{1}} \\
& =\frac{\kappa_{1}}{\tau_{2}} \frac{M_{1}}{M_{2}}+\frac{\tau_{2}}{\kappa_{1}} \frac{M_{2}}{M_{1}}=\text { constant }, \quad \frac{\kappa_{1}}{\tau_{2}}, \frac{M_{1}}{M_{2}}=\text { constant }
\end{aligned}
$$

and this is the intrinsic equation of a spacelike Bertrand curve.

## References

[1] Balgetir H., Bektas M., Inoguchi J, Null bertrand curves in Minkowski 3-space and their characterizations, Note di Matematica, 23, No:1, 713, 2004.
[2] Balgetir H., Bektas,M., Ergut,M., Bertrand curves for nonnull curves in 3-dimensional Lorentzian space, Hadronic Journal, 27, 229236, 2004.
[3] Bioche C., Sur les courbes de M. Bertrand, Bull. Soc.Math. France, 17, pp. 109-112, 1889.
[4] Bukcu B., Karacan M.K., On the involute and evolute curves of the spacelike curve with a spacelike binormal in Minkowski 3-Space, Int. J. Contemp. Math. Sciences, Vol. 2, no. 5, 221-232, 2007.
[5] Burke J.F., Bertrands curves associated with a pair of curves, Mathematics Magazine, Vol. 34, No. 1. pp.60-62, (1960).
[6] Ekmekci, N. ve Ilarslan K., On Bertrand curves and their characterization, Differential Geometry Dynamical Systems, 3, No.2, 17-24,201.
[7] Karacan M.K., Yayli Y., On the geodesics of tubular surfaces in Minkowski 3-space, Bulletin of the Malaysian Mathematical Sciences Society, 31 (1). pp. 1-10, 2008.
[8] Mellish A.P., Notes on differential geometry, Annals of Mathematics, Second Series, Vol. 32, No. 1, pp. 181-190, Jan. 1931.
[9] Oprea J., Differential Geometry and its Applications, Prentica Hall, Inc, 1997.
[10] Tajima J., On Bertrand curves, Tohoku Math. J., 18, pp.128-133,1920.
[11] Walrave J., Curves and Surfaces in Minkowski Space, PhD Thesis,1995.

# Respectable Graphs 

A.Satyanarayana Reddy<br>(Department of Mathematics and Statistics, Indian Institute of Technology Kanpur, India)<br>E-mail: satya@iitk.ac.in


#### Abstract

Two graphs $X$ and $Y$ are said to be respectable to each other if $\mathcal{A}(X)=$ $\mathcal{A}(Y)$. In this study we explore some graph theoretic and algebraic properties shared by the respectable graphs.


Key Words: Adjacency algebra, coherent algebra, walk regular graphs, vertex transitive graphs.

AMS(2010): 05C50, 05E40

## §1. Introduction

Let $A(X)$ (or simply $A$, if $X$ is clear from the context) be the adjacency matrix of a graph $X$. The set of all polynomials in $A$ with coefficients from the field of complex numbers $\mathbb{C}$ forms an algebra called the adjacency algebra of $X$, denoted by $\mathcal{A}(X)$. Let $\operatorname{dim}(\mathcal{A}(X))$ denote the dimension of $\mathcal{A}(X)$ as a vector space over $\mathbb{C}$. It is easy to see that $\operatorname{dim}(\mathcal{A}(X))$ is equal to degree of the minimal polynomial of $A$. Since $\operatorname{dim}(\mathcal{A}(X))$ is also equal $|\operatorname{spec}(A)|$ where $\operatorname{spec}(A)$ denote the set of all distinct eigenvalues of $A$ and $|B|$ denote the cardinality of the set $B$.

Definition 1 Two graphs $X$ and $Y$ are said to be respectable to each other if $\mathcal{A}(X)=\mathcal{A}(Y)$. In this case we say that either $X$ respects $Y$ or $Y$ respects $X$.

A graph $Y$ is said to be a polynomial in a graph $X$ if $A(Y) \in \mathcal{A}(X)$. For example, $K_{n}$ the complete graph is a polynomial in every connected regular graph with $n$ vertices. By definition if $X$ respects $Y$, then $X$ is a polynomial in $Y$ and $Y$ is a polynomial in $X$. In this study we explore some graph theoretic and algebraic properties shared by respectable graphs. In the remaining part of this section we will give some preliminaries required for this paper.

For two vertices $u$ and $v$ of a connected graph $X$, let $d(u, v)$ denote the length of the shortest path from $u$ to $v$. Then the diameter of a connected graph $X=(V, E)$ is $\max \{d(u, v)$ : $u, v \in V\}$. It is shown in Biggs [3] that if $X$ is a connected graph with diameter $\ell$, then $\ell+1 \leq \operatorname{dim}(\mathcal{A}(X)) \leq n$.

A graph $X_{1}=\left(V\left(X_{1}\right), E\left(X_{1}\right)\right)$ is said to be isomorphic to a graph $X_{2}=\left(V\left(X_{2}\right), E\left(X_{2}\right)\right)$, written $X_{1} \cong X_{2}$, if there is a one-to-one correspondence $\rho: V\left(X_{1}\right) \rightarrow V\left(X_{2}\right)$ such that $\left\{v_{1}, v_{2}\right\} \in E\left(X_{1}\right)$ if and only if $\left\{\rho\left(v_{1}\right), \rho\left(v_{2}\right)\right\} \in E\left(X_{2}\right)$. In such a case, $\rho$ is called an isomor-

[^9]phism of $X_{1}$ and $X_{2}$. An isomorphism of a graph $X$ onto itself is called an automorphism. The collection of all automorphisms of a graph $X$ is denoted by $\operatorname{Aut}(X)$. It is well known that $\operatorname{Aut}(X)$ is a group under composition of two maps. It is easy to see that $\operatorname{Aut}(X)=\operatorname{Aut}\left(X^{c}\right)$, where $X^{c}$ is the complement of the graph X. If $X$ is a graph with $n$ vertices we can think Aut $(X)$ as a subgroup of $S_{n}$. Under this correspondence, if a graph $X$ has $n$ vertices then $\operatorname{Aut}(X)$ consists of $n \times n$ permutation matrices and for each $g \in \operatorname{Aut}(X), P_{g}$ will denote the corresponding permutation matrix.

## §2. Graph Theoretic Properties

In this section we will see some graph theoretical properties shared by the respectable graphs. The next result gives a method to check whether a given permutation matrix is an element of $\operatorname{Aut}(X)$ or not.

Lemma 2.1(Biggs [3]) Let $A$ be the adjacency matrix of a graph $X$. Then $g \in A u t(X)$ is an automorphism of $X$ if and only if $P_{g} A=A P_{g}$.

The following result is immediate from the above lemma, also given by Paul.M. Weichsel [7].

Corollary 2.2 Let $X$ be a graph and $p(x)$ be a polynomial such that $p(X)$ is a graph. Then $\operatorname{Aut}(X) \subseteq \operatorname{Aut}(p(X))$.

Corollary 2.3 If the graph $X$ respects the graph $Y$, then $\operatorname{Aut}(X)=A u t(Y)$.
Lemma 2.4(Biggs [3]) A graph $X$ is regular if and only if $A(X) J=J A(X)$, where $J$ is a matrix with each entry is 1 .

The following result shows that any graph which is a polynomial in a regular graph is regular.

Corollary 2.5 Let $X$ be a regular graph. The any graph which is a polynomial in $X$ is also regular. In particular if $X$ respects $Y$, then $Y$ is regular.

Lemma 2.6(Biggs [3]) A graph $X$ is connected regular if and only if $J \in \mathcal{A}(X)$.
Corollary 2.7 If $X$ is a regular graph then $J$ is polynomial in either $A$ or $A^{c}$.
Proof For every graph $X$, either $X$ or $X^{c}$ is connected. Hence the result follows from the above lemma.

Corollary 2.8 Let $X$ be a connected regular graph, then $X^{c}$ is connected if and only if $X$ respects $X^{c}$.

Proof It is easy to verify that $X^{c}$ is also regular. Since $X$ is connected regular graph from Lemma 2.6 we have $J \in \mathcal{A}(X) \Rightarrow A\left(X^{c}\right)=J-I-A \in \mathcal{A}(X) \Rightarrow \mathcal{A}\left(X^{c}\right) \subseteq \mathcal{A}(X)$. Now it is
sufficient to prove that $X^{c}$ is connected if and only if $\mathcal{A}(X) \subseteq \mathcal{A}\left(X^{c}\right)$.
$X^{c}$ is connected $\Leftrightarrow J \in \mathcal{A}\left(X^{c}\right) \Leftrightarrow A \in \mathcal{A}\left(X^{c}\right) \Leftrightarrow \mathcal{A}(X) \subseteq \mathcal{A}\left(X^{c}\right)$.
Corollary 2.9 Let $X$ be a connected regular graph. If $X$ respects $Y$, then $Y$ is connected regular graph.

We say that a graph $X$ is walk-regular if, for each $s$, the number of closed walks of length $s$ starting at a vertex $v$ is independent of the choice of $v$.
Theorem 2.10([6]) Let $A$ be the adjacency matrix of a graph $X$. Then $X$ is walk-regular if and only if the diagonal entries of $A^{s} \forall s$ are all equal.

Corollary 2.11 Let $X$ be a walk regular graph and $p(x)$ be a polynomial such that $p(X)$ is a graph. Then $p(X)$ is walk regular.

Proof Let $A$ be the adjacency matrix of $X$. From the above theorem the diagonal entries of $A^{s} \forall s$ are all equal, so as for every element in $\mathcal{A}(X)$. As one of the basis for $\mathcal{A}(X)$ is $\left\{I, A, A^{2}, \ldots A^{l-1}\right\}$ where $l$ is the degree of the minimal polynomial of $A$.

From the above result we deduce that if $X$ be a walk regular graph and $X$ respects $Y$, then $Y$ is also walk regular graph.

Now we will see some symmetrical properties shared by the respectable graphs.

Definition 2 A graph $X=(V, E)$ is said to be vertex transitive if its automorphism group acts transitively on $V$. That is for any two vertices $x, y \in V, \exists g \in G$ such that $g(x)=y$.

Definition 3 A graph $X=(V, E)$ is said to be generously transitive if its automorphism group acts generously transitively on $V(X)$, i.e., if any $x, y \in V$ then $\exists g \in \operatorname{Aut}(X)$ such that $g(x)=y$ and $g(y)=x$.

Every generously transitive graph is transitive. From the Corollary 2.2 we have the following result.

Lemma 2.12 If $X$ is a generously transitive (or vertex transitive) graph and $Y$ is a polynomial in $X$, then $Y$ is also a generously transitive (or vertex transitive) graph.

## §3. Algebraic Properties

Let $X$ be a graph with $n$ vertices and $A$ be the adjacency matrix of $X$. By graph algebra of $X$, we mean a matrix subalgebra of $M_{n}(\mathbb{C})$ which contains $A$. For example $M_{n}(\mathbb{C})$ and $\mathcal{A}(X)$ are graph algebras of $X$. If the graph $X$ respects $Y$, then in this section we will show that the following three graph algebras of $X$ and $Y$ will coincide.

- The commutant algebra of a graph $Z$ is the set all matrices over $\mathbb{C}$ which commutes with adjacency matrix of $Z$.
- The coherent closure of a graph $Z$ is the smallest coherent algebra containing the adjacency matrix of $Z$.
- The centralizer algebra of a graph $Z$ is the set all matrices which commute with all automorphisms of $Z$.


### 3.1 Coherent Closure of a Graph

Definition 4 Hadamard product of two $n \times n$ matrices $A$ and $B$ is denoted by $A \odot B$ and is defined as $(A \odot B)_{x y}=A_{x y} B_{x y}$.

Definition 5 Two $n \times n$ matrices $A$ and $B$ are said to be disjoint if their Hadamard product is the zero matrix.

Definition $6 \quad A$ sub algebra of $M_{n}(\mathbb{C})$ is called coherent if it contains the matrices I and $J$ and if it is closed under conjugate-transposition and Hadamard multiplication.

The following result is well known.

Theorem 3.1 Every coherent algebra contains unique basis of disjoint 0-1 matrices.
We call the unique basis containing disjoint $0-1$ matrices as a standard basis.
Corollary 3.2 Every 0,1-matrix in a coherent algebra is sum of one or more matrices in its standard basis.

Proof Let $\mathcal{M}$ be a coherent algebra over $\mathbb{C}$ with standard basis $\left\{M_{1}, \ldots M_{t}\right\}$. Let $B \in \mathcal{M}$ be a 0,1 -matrix. Then $B=\sum_{i=1}^{t} a_{i} M_{i}$ where $a_{i} \in \mathbb{C} . B=B \odot B=\sum_{i=1}^{t} a_{i}^{2} M_{i} \Rightarrow a_{i}^{2}=a_{i}$. Hence the result follows.

Observation 3.3 The intersection of coherent algebras is again a coherent algebra.
Definition 7 Let $X=(V, E)$ be a graph with adjacency matrix $A$ then any coherent algebra which contains $A$ is called coherent algebra of $X$.

Definition 8 If $X=(V, E)$ be a graph and $A$ is its adjacency matrix then coherent closure of $X$, denoted by $\langle\langle A\rangle\rangle$ or $\mathcal{C C}(X)$, is the smallest coherent algebra containing $A$.

Since $A\left(X^{c}\right)=J-I-A(X)$ consequently $A(X), A\left(X^{c}\right) \in \mathcal{C C}(X) \cap \mathcal{C C}\left(X^{c}\right)$, hence we get the following lemma.

Lemma 3.4 For every graph $X, \mathcal{C C}(X)=\mathcal{C C}\left(X^{c}\right)$.

Lemma 3.5 If the graph $X$ respects $Y$, then $\mathcal{C C}(X)=\mathcal{C C}(Y)$.
Proof Since $X$ respects $Y$, we have $\mathcal{A}(X)=\mathcal{A}(Y) \subseteq \mathcal{C C}(Y)$. Consequently $\mathcal{C C}(Y)$ is a coherent algebra containing $A(X)$ but by definition $\mathcal{C C}(X)$ is the smallest coherent algebra containing $A(X)$. So $\mathcal{C C}(X) \subseteq \mathcal{C C}(Y)$. Similarly we can prove $\mathcal{C C}(Y) \subseteq \mathcal{C C}(X)$. Hence the result follows.

Clearly, the converse of this result is not true as $\mathcal{C C}(X)=\mathcal{C C}\left(X^{c}\right)$, but $X$ need not respect $X^{c}$.

### 3.2 Centralizer Algebra of a Graph

Definition 9 Let $G$ be a subset of $n \times n$ permutation matrices forming a group. Then $\mathcal{V}_{\mathbb{C}}(G)=$ $\left\{A \in M_{n}(\mathbb{C}): P A=A P \forall P \in G\right\}$ forms an algebra over $\mathbb{C}$ called the centralizer algebra of the group $G$.

Definition 10 If $G$ is a group acting on a set $V$, then $G$ also acts on $V \times V$ by $g(x, y)=$ $(g(x), g(y))$. The orbits of $G$ on $V \times V$ are called orbitals. In the context of graphs, the orbitals of graph $X$ are orbitals of its automorphism group $\operatorname{Aut}(X)$ acting on the vertex set of $X$. That is, the orbitals are the orbits of the arcs/non-arcs of the graph $X=(V, E)$. The number of orbitals is called the rank of $X$.

An orbital can be represented by a $0-1$ matrix $M$ where $M_{i j}$ is 1 if $(i, j)$ belongs to the orbital. We can associate directed graphs to these matrices. If the matrices are symmetric, then these can be treated as undirected graphs.

## Observation 3.6

- The ' 1 ' entries of any orbital matrix are either all on the diagonal or all are off diagonal.
- The orbitals containing 1's on the diagonal will be called diagonal orbitals.

Definition 11 The centralizer algebra of a graph $X$ denoted by $\mathcal{V}(X)$ is the centralizer algebra of its automorphism group.

Theorem 3.7([4]) $\quad \mathcal{V}_{\mathbb{C}}(G)$ is a coherent algebra and orbitals of Aut $X$ acting on the vertex set of $X$ form its unique 0-1 matrix basis.
$\mathcal{V}(X)=\mathcal{V}\left(X^{c}\right)$ follows from the fact that $\operatorname{Aut}(X)=\operatorname{Aut}\left(X^{c}\right) \cdot \mathcal{C C}(X)$ is the smallest coherent algebra of $X$ and $\mathcal{V}(X)$ is a coherent algebra of $X$ so $\mathcal{C C}(X) \subseteq \mathcal{V}(X)$. So for any graph $X$ we have $\mathcal{A}(X) \subseteq \mathcal{C C}(X) \subseteq \mathcal{V}(X)$. The following result follows from the Corollary 2.3.

Lemma 3.8 If the graph $X$ respects the graph $Y$, then $\mathcal{V}(X)=\mathcal{V}(Y)$.
Now we will see a consequence of above result. For that we need the following definition.
Definition 12(Robert A.Beezer [1]) A graph $X=(V, E)$ is orbit polynomial graph if each orbital matrix is member of $\mathcal{A}(X)$. That is each orbital matrix is a polynomial in $A$.

Lemma $3.9 X$ is an orbit polynomial graph if and only if $\mathcal{A}(X)=\mathcal{V}(X)$.
If $X$ is an orbit polynomial graph, then we have $\mathcal{A}(X)=\mathcal{C C}(X)=\mathcal{V}(X)$.
Corollary 3.10 Let $X$ be an orbit polynomial graph and suppose $X$ respects the graph $Y$, then $Y$ is also an orbit polynomial graph.

Corollary 3.11 If $X$ is an orbit polynomial graph and $X^{c}$ is connected then $X^{c}$ is orbit polynomial graph.

If $X$ is an orbit polynomial graph and $X^{c}$ is connected, then we have $\mathcal{A}(X)=\mathcal{A}\left(X^{c}\right)=$ $\mathcal{C C}(X)=\mathcal{C C}\left(X^{c}\right)=\mathcal{V}(X)=V\left(X^{c}\right)$.

### 3.3 Commutant algebra of a graph

The commutant algebra of graph $X$, denoted by $C[X]$ is the set of all matrices which commutes with $A$. It is shown in (Davis [5]) that $\operatorname{dim}(C[A])=$ sum of the squares of the multiplicities of eigenvalues of $A$. Hence the following lemma.

Lemma $3.12 \mathcal{A}(X)=C[X]$ if and only if all eigenvalues of $X$ are distinct.
Lemma 3.13 If the graph $X$ respects the graph $Y$ then $C[X]=C[Y]$.
Proof Notice that

$$
B \in C[X] \Leftrightarrow B A(X)=A(X) B \Leftrightarrow B A(Y)=A(Y) B \Leftrightarrow B \in C[Y]
$$

We get the result.

## §4 Polynomial Equivalence

Let $\mathcal{G}_{n}$ be the set of all graphs with $n$ vertices. We define a relation $\mathcal{R}$ on $\mathcal{G}_{n}$ as $X \mathcal{R} Y \Leftrightarrow$ $X$ respects $Y$. It is easy to see that $\mathcal{R}$ is an equivalence relation on $\mathcal{G}_{n}$. Now for a given graph $X$, our objective is to find the equivalence class $[X]$ under the equivalence relation $\mathcal{R}$. First we identify a set $[X]$ with a set in polynomial algebra $\mathbb{C}[x]$. For that we need the following notations and definitions.
$\mathbb{C}[A]$ denote the set of all matrices which are polynomials in the square matrix $A$. It is easy to see that $\mathbb{C}[A] \cong \mathbb{C}[x] /\langle p(x)\rangle$ where $\langle p(x)\rangle$ is the ideal in $\mathbb{C}[x]$ generated by $p(x)$, which is the minimal polynomial of $A$. Consequently if $B \in \mathbb{C}[A]$, then there exists a unique polynomial $f_{B}(x)$ called representor polynomial of $B$ such that $\operatorname{deg}\left(f_{B}(X)\right) \leq \operatorname{deg}(p(x))$ and $f_{B}(A)=B$.

Definition 13 Let $A$ be a square matrix and $f(x)$ be a polynomial. We say that $f(x)$ respects $\operatorname{spec}(A)$ if $f\left(\lambda_{i}\right) \neq f\left(\lambda_{j}\right)$ for $\lambda_{i}$ and $\lambda_{j}$ distinct eigenvalues of $A$.

The following result is given by Paul M.Weichsel [7].

Lemma 4.1 Let $A$ be diagonalizable matrix over a field and $f(x) \in \mathbb{C}[x]$. Then $f(x)$ respects $\operatorname{spec}(A)$ if and only if there exists a polynomial $g(x) \in \mathbb{C}[x]$ such that $g(f(A))=A$.

Proof Let $B=f(A)$. Clearly $\mathbb{C}[B] \subseteq \mathbb{C}[A]$. Since $A$ is diagonalizable, so is $B$. Consequently $A$ is a polynomial in $B$ if and only if $\mathbb{C}[B]=\mathbb{C}[A]$ if and only if $|\operatorname{spec}(A)|=|\operatorname{spec}(B)|$ if and only if $f(x)$ respects $\operatorname{spec}(A)$.

Now let $A$ be the adjacency matrix of a graph $X$ and we denote

$$
\begin{aligned}
& F_{X}=\{f(x) \in \mathbb{C}[x] \mid \operatorname{deg}(f(x)) \leq \operatorname{deg}(p(x)) \text { and } f(A) \text { is a } 0,1 \text {-matrix }\}, \\
& H_{X}=\left\{g(x) \in F_{X} \mid g(x) \text { respects } \operatorname{spec}(A)\right\}
\end{aligned}
$$

Now one can easily verify that finding the set $[X]$ is equivalent to finding the set $H_{X}$. By definition, in order to find $H_{X}$ we need to find $F_{X}$ but for a given graph $X$ finding $F_{X}$ seems difficult. Let $X$ be a graph with the property $\mathcal{A}(X)=\mathcal{C C}(X)$, then from Corollary 3.2 it is easy to evaluate $F_{X}$. Distance regular graphs and orbit polynomial graphs satisfy $\mathcal{A}(X)=\mathcal{C C}(X)$ for details one can refer Robert A.Beezer [2] and Paul M.Weichsel [7].

The following theorem shows that if $X$ is a connected vertex transitive graph with a prime number of vertices then $X$ respects $Y$ if and only if $\operatorname{Aut}(X)=\operatorname{Aut}(Y)$.

Theorem 4.2(Robert A.Beezer [2]) Suppose that $X$ is a connected, vertex transitive graph with a prime number of vertices. Let $p(x)$ be a polynomial such that $p(X)$ is a connected graph, and $\operatorname{Aut}(X)=\operatorname{Aut}(p(X))$. Then $p(x)$ respects spec $((A(X))$.

Comments In spite of these results, there are many properties which are not shared by the respectable graphs. We illustrate few of them with examples. Let $C_{n}$ denote the cycle graph with $n$ vertices, then $C_{n}$ respects $C_{n}^{c}$ for $n \geq 5$. It is known that $C_{2 n}$ is bipartite for every $n$, but $C_{2 n}^{c}$ is not bipartite for $n \geq 3$. For $n \geq 3, C_{2 n}$ is Eulerain graph but $C_{2 n}^{c}$ is not. $C_{n}$ is planar graph $\forall n$ where as $C_{n}^{c}$ is not planar for $n \geq 9$ as every finite, simple, planar graph has a vertex of degree less than 6. Petersen graph is not Hamiltonian graph but from Dirac's theorem its compliment (respects Petersen graph) is a Hamiltonian graph.

## references

[1] Robert A.Beezer, Trivalent orbit polynomial graphs, Linear Algebra and its Applications, Volume 73, January 1986, Pages 133-146.
[2] Robert A.Beezer, A disrespectful polynomial, Linear Algebra and its Applications, Volume 128, January 1990, Pages 139-146
[3] N.L.Biggs, Algebraic Graph Theory (second ed.), Cambridge University Press, Cambridge (1993).
[4] A.E.Brouwer, A. M. Cohen, A.Neumaier, Distance Regular Graphs, Springer-Verlag, 1989.
[5] Philip J. Davis, Circulant Matrices, A Wiley-interscience publications, 1979.
[6] C. D.Godsil and B. D. cKay, Feasibility conditions for the existence of walk-regular graphs, Linear Algebra and its Applications, 30:51-61(1980).
[7] A.Satyanarayana Reddy and Shashank K Mehta, Pattern polynomial graphs, Contributed to Mathematics ArXiv via http://arxiv.org/abs/1106.4745.
[8] Paul M.Weichsel, Polynomials on graphs, Linear Algebra and its Applications, 93:179-186 (1987).

# Common Fixed Points for Pairs of Weakly Compatible Mappings 

Rakesh Tiwari<br>(Department of Mathematics, Govt. Arts and Science College, Durg (C.G.), 491001, India)<br>S.K.Shrivastava<br>(Department of Mathematics, Deen Dayal Upadhyay University, Gorakhpur (U. P.), 273009, India)<br>E-mail: rakeshtiwari66@gmail.com, sudhirpr66@rediffmail.com


#### Abstract

In this note we establish a common fixed point theorem for a quadruple of self mappings satisfying a common (E.A) property on a metric space satisfying weakly compatibility and a generalized $\Phi$ - contraction. Our results improve and extend some known results.


Key Words: Common fixed points, weakly compatible mappings, generalized $\Phi$ - contraction, a common (E.A) property, Smarandache metric multi-space.

AMS(2010): 47H10, 54 H 25

## §1. Introduction

For an integer $n \geq 1$, a Smarandache metric multi-space $\widetilde{S}$ is a union $\bigcup_{i=1}^{n} A_{i}$ of spaces $A_{1}, A_{2}, \cdots$, $A_{n}$, distinct two by two with metrics $\rho_{1}, \rho_{2}, \cdots, \rho_{n}$ such that $\left(A_{i}, \rho_{i}\right)$ is a metric space for integers $1 \leq i \leq n$. In 1986, the notion of compatible mappings which generalized commuting mappings, was introduced by Jungck [3]. This has proven useful for generalization of results in metric fixed point theory for single-valued as well as multi-valued mappings. Further in 1998, the more general class of mappings called weakly compatible mappings was introduced by Jungck and Rhoades [4]. Recall that self mappings S and T of a metric space $(X, d)$ are called weakly compatible if $S x=T x$ for some $x \in X$ implies that $S T x=T S x$.

Recently Aamri et al. [1] introduced the following notion for a pair of maps as:

Definition 1.1 Let $S$ and $T$ be two self mappings of a metric space $(X, d) . S$ and $T$ are said to satisfy the property (E.A), if there exists a sequence $\left\{x_{n}\right\}$ in $X$ such that $\lim _{n \rightarrow \infty} T x_{n}=$ $\lim _{n \rightarrow \infty} S x_{n}=t$, for some $t \in X$.

Most recently, Y. Liu et al. [5] defined a common property (E.A) for pairs of mappings as

[^10]follows:

Definition 1.2 Let $A, B, S, T: X \rightarrow X$. The pairs $(A, S)$ and $(B, T)$ satisfy a common property (E.A) if there exist two sequences $\left\{x_{n}\right\}$ and $\left\{y_{n}\right\}$ such that

$$
\lim _{n \rightarrow \infty} A x_{n}=\lim _{n \rightarrow \infty} S x_{n}=\lim _{n \rightarrow \infty} B y_{n}=\lim _{n \rightarrow \infty} T y_{n}=t \in X
$$

If $B=A$ and $S=T$ in above, we obtain the definition of property (E.A).
Example 1.3 Let $A, B, S$ and $T$ be self maps on $X=[0,1]$, with the usual metric $d(x, y)=$ $|x-y|$, defined by:

$$
\begin{aligned}
& A x=\left\{\begin{array}{ccc}
1-\frac{x}{2} & \text { when } & x \in\left[0, \frac{1}{2}\right), \\
1 & \text { when } & \left.x \in\left[\frac{1}{2}, 1\right]\right) .
\end{array}\right. \\
& S x=\left\{\begin{array}{ccc}
1-2 x & \text { when } & x \in\left[0, \frac{1}{2}\right), \\
1 & \text { when } & \left.x \in\left[\frac{1}{2}, 1\right]\right) .
\end{array}\right.
\end{aligned}
$$

$B x=1-x$ and $T x=1-\frac{x}{3}, \forall x \in X$. Let $\left\{x_{n}\right\}$ and $\left\{y_{n}\right\}$ be a sequences defined by $x_{n}=\frac{1}{n+1}$ and $y_{n}=\frac{1}{n^{2}+1}$, then $\lim _{n \rightarrow \infty} A x_{n}=\lim _{n \rightarrow \infty} S x_{n}=\lim _{n \rightarrow \infty} B y_{n}=\lim _{n \rightarrow \infty} T y_{n}=1 \in X$. Thus a common (E.A) property is satisfied.

In this paper we prove some common fixed point theorems for a quadruple of weak compatible self mappings of a metric space satisfying a common (E.A) property, a special Smarandache metric multi-space $\bigcup_{i=1}^{n}\left(A_{i}, \rho_{i}\right)$ for $n=1$ and a generalized $\Phi$-contraction. These theorems extend and generalize results of Pathak et al. [6] and [7].

## §2. Preliminaries

Now onwards, we denote by $\Phi$ the collection of all functions $\varphi:[0, \infty) \rightarrow[0, \infty)$ which are upper semi-continuous from the right, non-decreasing and satisfy $\lim _{s \rightarrow t+} \sup \varphi(s)<t, \varphi(t)<t$ for all $t>0$.

Let $X$ denote a metric space endowed with metric $d$ and let $\mathbb{N}$ denote the set of natural numbers.

Now, let $A, B, S$ and $T$ be self-mappings of $X$ such that

$$
\begin{align*}
& {\left[d^{p}(A x, B y)+a d^{p}(S x, T y)\right] d^{p}(A x, B y)} \\
& \leq a \max \left\{d^{p}(A x, S x) d^{p}(B y, T y), d^{q}(A x, T y) d^{q^{\prime}}(B y, S x)\right\} \\
& \quad+\max \left\{\varphi_{1}\left(d^{2 p}(S x, T y)\right), \varphi_{2}\left(d^{r}(A x, S x) d^{r^{\prime}}(B y, T y)\right),\right. \\
& \varphi_{3}\left(d^{s}(A x, T y) d^{s^{\prime}}(B y, S x)\right), \\
& \varphi_{4}\left(\frac{1}{2}\left[d^{l}(A x, T y) d^{l^{\prime}}(A x, S x)+d^{l}(B y, S x)\right) d^{l^{\prime}}(B y, T y)\right\} \tag{2.1}
\end{align*}
$$

for all $x, y \in X, \varphi_{i} \in \Phi(i=1,2,3,4), a, p, q, q^{\prime}, r, r^{\prime}, s, s^{\prime}, l, l^{\prime} \geq 0$ and $2 p=q+q^{\prime}=r+r^{\prime}=$ $s+s^{\prime}=l+l^{\prime}$. The condition (2.1) is commonly called a generalized $\Phi$-contraction.

## §3. Main Results

The following theorems are our main results in this section.

Theorem 3.1 Let $A, B, S$ and $T$ be self mappings of a metric space ( $X, d$ ) satisfying (2.1).If the pairs $(A, S)$ and $(B, T)$ satisfy a common (E.A) property, are weakly compatible and that $T(X)$ and $S(X)$ are closed subsets of $X$, then $A, B, S$ and $T$ have a unique common fixed point in $X$.

Proof. Since $(A, S)$ and $(B, T)$ satisfy a common property (E.A). Then there exist two sequences $\left\{x_{n}\right\}$ and $\left\{y_{n}\right\}$ such that

$$
\lim _{n \rightarrow \infty} A x_{n}=\lim _{n \rightarrow \infty} S x_{n}=\lim _{n \rightarrow \infty} B y_{n}=\lim _{n \rightarrow \infty} T y_{n}=z
$$

for some $z \in X$. Assume that $S(X)$ and $T(X)$ are closed subspaces of $X$. Then, $z=S u=T v$ for some $u, v \in X$. Then by using (2.1) with $x=x_{n}$ and $y=v$, we have

$$
\begin{aligned}
& {\left[d^{p}\left(A x_{n}, B v\right)+a d^{p}\left(S x_{n}, T v\right)\right] d^{p}\left(A x_{n}, B v\right) \leq a \max \left\{d^{p}\left(A x_{n}, S x_{n}\right) d^{p}(B v, T v),\right.} \\
& \left.d^{q}\left(A x_{n}, T v\right) d^{q^{\prime}}\left(B v, S x_{n}\right)\right\}+\max \left\{\varphi_{1}\left(d^{2 p}\left(S x_{n}, T v\right)\right),\right. \\
& \varphi_{2}\left(d^{r}\left(A x_{n}, S x_{n}\right) d^{r^{\prime}}(B v, T v)\right), \varphi_{3}\left(d^{s}\left(A x_{n}, T v\right) d^{s^{\prime}}\left(B v, S x_{n}\right)\right), \\
& \left.\varphi_{4}\left(\frac{1}{2}\left[d^{l}\left(A x_{n}, T v\right) d^{l^{\prime}}\left(A x_{n}, S x_{n}\right)+d^{l}\left(B v, S x_{n}\right)\right) d^{l^{\prime}}(B v, T v)\right]\right),
\end{aligned}
$$

taking $\lim _{n \rightarrow \infty}$, we obtain

$$
\begin{aligned}
{\left[d^{p}(z, B v)+a d^{p}(z, T v)\right] d^{p}(z, B v) \leq } & a \max \left\{d^{p}(z, z) d^{p}(B v, z), d^{q}(z, T v) d^{q^{\prime}}(B v, z)\right\} \\
& +\max \left\{\varphi_{1}\left(d^{2 p}(z, T v)\right), \varphi_{2}\left(d^{r}(z, z) d^{r^{\prime}}(B v, z)\right),\right. \\
& \varphi_{3}\left(d^{s}(z, T v) d^{s^{\prime}}(B v, z)\right), \varphi_{4}\left(\frac { 1 } { 2 } \left[d^{l}(z, T v) d^{l^{\prime}}(z, z)\right.\right. \\
& \left.\left.\left.\left.+d^{l}(B v, z)\right) d^{l^{\prime}}(B v, z)\right]\right)\right\},
\end{aligned}
$$

or

$$
\begin{aligned}
d^{2 p}(z, B v) \leq & \max \left\{\varphi_{1}(0), \varphi_{2}(0), \varphi_{3}(0), \varphi_{4}\left(\frac{1}{2} d^{l+l^{\prime}}(B v, z)\right)\right\}, \\
d^{2 p}(z, B v) \leq & \max \left\{\varphi_{1}\left(d^{2 p}(z, B v)\right), \varphi_{2}\left(d^{r+r^{\prime}}(z, B v)\right.\right. \\
& \left.\varphi_{3}\left(d^{s+s^{\prime}}(z, B v)\right), \varphi_{4}\left(\frac{1}{2} d^{l+l^{\prime}}(B v, z)\right)\right\}
\end{aligned}
$$

This together with a well known result of Chang [2] which states that if $\varphi_{i} \in \Phi$ where $i \in I$ (some indexing set), then there exists a $\varphi \in \Phi$ such that $\max \left\{\varphi_{i}, i \in I\right\} \leq \varphi(t)$ for all $t>0$; imply

$$
d^{2 p}(z, B v) \leq \varphi\left(d^{2 p}(z, B v)\right)<d^{2 p}(z, B v)
$$

a contradiction. This implies that $z=B v$. Therefore $T v=z=B v$. Hence it follows by the weak compatibility of the pair $(B, T)$ that $B T v=T B v$, that is $B z=T z$.

Now, we shall show that $z$ is a common fixed point of $B$ and $T$. For this put $x=x_{n}$ and $y=z$ in (2.1), we have

$$
\begin{aligned}
& {\left[d^{p}\left(A x_{n}, B z\right)+a d^{p}\left(S x_{n}, T z\right)\right] d^{p}\left(A x_{n}, B z\right) \leq a \max \left\{d^{p}\left(A x_{n}, S x_{n}\right) d^{p}(B z, T z),\right.} \\
& \\
& \left.\qquad d^{q}\left(A x_{n}, T z\right) d^{q^{\prime}}\left(B z, S x_{n}\right)\right\}+\max \left\{\varphi_{1}\left(d^{2 p}\left(S x_{n}, T z\right)\right),\right. \\
& \varphi_{2}\left(d^{r}\left(A x_{n}, S x_{n}\right) d^{r^{\prime}}(B z, T z)\right), \varphi_{3}\left(d^{s}\left(A x_{n}, T z\right) d^{s^{\prime}}\left(B z, S x_{n}\right)\right), \\
& \varphi_{4}\left(\frac{1}{2}\left[d^{l}\left(A x_{n}, T z\right) d^{l^{\prime}}\left(A x_{n}, S x_{n}\right)+d^{l}\left(B z, S x_{n}\right)\right) d^{l^{\prime}}(B z, T z)\right\}
\end{aligned}
$$

Letting $n \rightarrow \infty$ with the help of the fact that $\lim _{n \rightarrow \infty} A x_{n}=z=\lim _{n \rightarrow \infty} S x_{n}$ and $B z=T z$, we get

$$
\begin{gathered}
{\left[d^{p}(z, B z)+a d^{p}(z, T z)\right] d^{p}(z, B z) \leq a \max \left\{d^{p}(z, z) d^{p}(B z, z), d^{q}(z, T z) d^{q^{\prime}}(B z, z)\right\}} \\
+\max \left\{\varphi_{1}\left(d^{2 p}(z, T z)\right), \varphi_{2}\left(d^{r}(z, z) d^{r^{\prime}}(B z, z)\right), \varphi_{3}\left(d^{s}(z, T z) d^{s^{\prime}}(B z, z)\right)\right. \\
\left.\left.\varphi_{4}\left(\frac{1}{2}\left[d^{l}(z, T z) d^{l^{\prime}}(z, z)+d^{l}(B z, z)\right) d^{l^{\prime}}(B z, z)\right]\right)\right\}
\end{gathered}
$$

or

$$
d^{2 p}(z, B z)+a d^{2 p}(z, B z) \leq a d^{q+q^{\prime}}(B z, z)+\max \left\{\varphi_{1}\left(d^{2 p}(z, B z)\right)\right.
$$

$$
\left.\varphi_{2}(0), \varphi_{3}\left(d^{s+s^{\prime}}(z, B z)\right), \varphi_{4}(0)\right\}
$$

or

$$
\begin{gathered}
\left.(1+a) d^{2 p}(z, B z) \leq a d^{q+q^{\prime}}(B z, z)\right\}+\max \left\{\varphi_{1}\left(d^{2 p}(z, B z)\right)\right. \\
\left.\varphi_{2}(0), \varphi_{3}\left(d^{s+s^{\prime}}(z, B z)\right), \varphi_{4}(0)\right\} \\
d^{2 p}(z, B z) \leq \frac{a}{1+a} d^{q+q^{\prime}}(B z, z)+\frac{1}{1+a} \max \left\{\varphi_{1}\left(d^{2 p}(z, B z)\right)\right. \\
\left.\varphi_{2}(0), \varphi_{3}\left(d^{s+s^{\prime}}(z, B z)\right), \varphi_{4}(0)\right\} \\
<d^{2 p}(z, B z)
\end{gathered}
$$

a contradiction. So $z=B z=T z$. Thus $z$ is a common fixed point of $B$ and $T$.
Similarly we can prove that $z$ is a common fixed point of $A$ and $S$. Thus $z$ is the common fixed point of $A, B, S$ and $T$. The uniqueness of $z$ as a common fixed point of $A, B, S$ and $T$ can easily be verified.

Remark 3.3 Our Theorem 3.1 extends theorem 2.1 of Pathak et al. [6].
In Theorem 3.1, if we put $a=0$ and $\varphi_{i}(t)=h t(i=1,2,3,4)$, where $0<h<1$, we get the following corollary:

Corollary 3.4 Let $A, B, S$ and $T$ be self mappings of a metric space $X$. If the pairs $(A, S)$ and $(B, T)$ satisfy a common (E.A) property and

$$
\begin{array}{r}
d^{2 p}(A x, B y) \leq h \max \left\{d^{2 p}(S x, T y), d^{r}(A x, S x) d^{r^{\prime}}(B y, T y), d^{s}(A x, T y)\right. \\
\left.\left.d^{s^{\prime}}(B y, S x)\right), \frac{1}{2}\left[d^{l}(A x, T y) d^{l^{\prime}}(A x, S x)+d^{l}(B y, S x)\right) d^{l^{\prime}}(B y, T y)\right\} \tag{2.2}
\end{array}
$$

for all $x, y \in X, \varphi_{i} \in \Phi(i=1,2,3,4)$, a, $p, q, q^{\prime}, r, r^{\prime}, s, s^{\prime}, l, l^{\prime} \geq 0$ and $2 p=q+q^{\prime}=r+r^{\prime}=$ $s+s^{\prime}=l+l^{\prime}$. If the pairs $(A, S)$ and $(B, T)$ are weakly compatible and that $T(X)$ and $S(X)$ are closed, then $A, B, S$ and $T$ have a unique common fixed point in $X$.

Especially when

$$
\begin{aligned}
& \max \left\{d^{2 p}(S x, T y), d^{r}(A x, S x) d^{r^{\prime}}(B y, T y), d^{s}(A x, T y) d^{s^{\prime}}(B y, S x)\right) \\
& \left.\frac{1}{2}\left[d^{l}(A x, T y) d^{l^{\prime}}(A x, S x)+d^{l}(B y, S x)\right) d^{l^{\prime}}(B y, T y)\right\}=d^{2 p}(S x, T y)
\end{aligned}
$$

it generalizes Corollary 3.9 of Pathak et al. [7].

In Theorem 3.1, if we take $S=T=I_{X}$ (the identity mapping on $X$ ), then we have the following corollary:

Corollary 3.5 Let $A$ and $B$ be self mappings of a complete metric space $X$ satisfying the following condition:

$$
\begin{aligned}
& {\left[d^{p}(A x, B y)+a d^{p}(x, y)\right] d^{p}(A x, B y) \leq a \max \left\{d^{p}(A x, x) d^{p}(B y, y)\right.} \\
& \left.\quad d^{q}(A x, y) d^{q^{\prime}}(B y, x)\right\}+\max \left\{\varphi_{1}\left(d^{2 p}(x, y)\right), \varphi_{2}\left(d^{r}(A x, x) d^{r^{\prime}}(B y, y)\right),\right. \\
& \quad \varphi_{3}\left(d^{s}(A x, y) d^{s^{\prime}}(B y, x)\right), \varphi_{4}\left(\frac{1}{2}\left[d^{l}(A x, y) d^{l^{\prime}}(A x, x)+d^{l}(B y, x)\right) l^{l^{\prime}}(B y, y)\right\}
\end{aligned}
$$

for all $x, y \in X, \varphi_{i} \in \Phi(i=1,2,3,4), a, p, q, q^{\prime}, r, r^{\prime}, s, s^{\prime}, l, l^{\prime} \geq 0$ and $2 p=q+q^{\prime}=r+r^{\prime}=$ $s+s^{\prime}=l+l^{\prime}$, then $A$ and $B$ have a unique common fixed point in $X$.

As an immediate consequences of Theorem 3.1 with $S=T$, we have the following:
Corollary 3.6 Let $A, B$, and $S$ be self-mappings of $X$ such that $(A, S)$ and $(B, S)$ satisfy a common (E.A) property and

$$
\begin{gather*}
d^{2 p}(A x, B y) \leq a \max \left\{d^{p}(A x, S x) d^{p}(B y, S y), d^{q}(A x, S y) d^{q^{\prime}}(B y, S x)\right\} \\
+\max \left\{\varphi_{2}\left(d^{r}(A x, S x) d^{r^{\prime}}(B y, S y)\right), \varphi_{3}\left(d^{s}(A x, S y) d^{s^{\prime}}(B y, S x)\right),\right. \\
\varphi_{4}\left(\frac{1}{2}\left[d^{l}(A x, S y) l^{l^{\prime}}(A x, S x)+d^{l}(B y, S x)\right) d^{l^{\prime}}(B y, S y)\right\} \tag{2.3}
\end{gather*}
$$

for all $x, y \in X, \varphi_{i} \in \Phi(i=1,2,3,4), a, p, q, q^{\prime}, r, r^{\prime}, s, s^{\prime}, l, l^{\prime} \geq 0$ and $2 p=q+q^{\prime}=r+r^{\prime}=$ $s+s^{\prime}=l+l^{\prime}$. If the pairs $(A, S)$ and $(B, S)$ are weakly compatible and that $S(X)$ is closed, then $A, B$ and $S$ have a unique common fixed point in $X$.

Theorem 3.7 Let $S, T$ and $A_{n}(n \in \mathbb{N})$ be self mappings of a metric space $(X, d)$. Suppose further that the pairs $\left(A_{2 n-1}, S\right)$ and $\left(A_{2 n}, T\right)$ are weakly compatible for any $n \in \mathbb{N}$ and satisfying a common (E.A) property. If $S(X)$ and $T(X)$ are closed and that for any $i \in N$, the following condition is satisfied for all $x, y \in X$

$$
\begin{aligned}
& {\left[d^{p}\left(A_{i} x, A_{i+1} y\right)+a d^{p}(S x, T y)\right] d^{p}\left(A_{i} x, A_{i+1} y\right)} \\
& \qquad a \max \left\{d^{p}\left(A_{i} x, S x\right) d^{p}\left(A_{i+1} y, T y\right),\right. \\
& \left.d^{q}\left(A_{i} x, T y\right) d^{q^{\prime}}\left(A_{i+1} y, S x\right)\right\}+\max \left\{\varphi_{1}\left(d^{2 p}(S x, T y)\right),\right. \\
& \varphi_{2}\left(d^{r}\left(A_{i} x, S x\right) d^{r^{\prime}}\left(A_{i+1} y, T y\right)\right), \varphi_{3}\left(d^{s}\left(A_{i} x, T y\right) d^{s^{\prime}}\left(A_{i+1} y, S x\right)\right), \\
& \varphi_{4}\left(\frac{1}{2}\left[d^{l}\left(A_{i} x, T y\right) d^{l^{\prime}}\left(A_{i} x, S x\right)+d^{l}\left(A_{i+1} y, S x\right)\right) d^{l^{\prime}}\left(A_{i+1} y, T y\right)\right\}
\end{aligned}
$$

where $\varphi_{i} \in \Phi(i=1,2,3,4), a, p, q, q^{\prime}, r, r^{\prime}, s, s^{\prime}, l, l, \geq 0$ and $2 p=q+q^{\prime}=r+r^{\prime}=s+s^{\prime}=l+l^{\prime}$, then $S, T$ and $A_{n}(n \in \mathbb{N})$ have a common fixed point in $X$.

Acknowledgement: The authors are very grateful to Prof. H. K. Pathak for his valuable suggestions regarding this paper.

## References

[1] M.Aamri and D. El Moutawakil, Some new common fixed point theorems under strict contractive conditions, J. Math. Anal. Appl. 270(2002), 181-188.
[2] S. S.Chang, A common fixed point theorem for commuting mappings, Math. Japon, 26 (1981), 121-129.
[3] G.Jungck, Compatible mappings and common fixed points, Int. J. Math. Math. Sci., $\mathbf{9}$ (1986), 771-779.
[4] G.Jungck and B.E.Rhoades, Fixed points for set valued functions without continuity, Indian J. Pure Appl. Math., 29 (3)(1998), 227-238.
[5] W.Liu, J.Wu, Z. Li, Common fixed points of single-valued and multi-valued maps, Int. J. Math. Math. Sc., 19(2005), 3045-3055.
[6] H.K.Pathak, S.N.Mishra and A.K.Kalinde, Common fixed point theorems with applications to non-linear integral equations, Demonstratio Math., XXXII(3) (1999), 547-564.
[7] H.K.Pathak, Y.J.Cho and S.M.Kang, Common fixed points of biased maps of type (A) and applications, Int.J. Math. and Math. Sci., 21(4) (1999), 681-694.

# Some Results on 

# Generalized Multi Poly-Bernoulli and Euler Polynomials 

Hassan Jolany, Hossein Mohebbi<br>(Department of Mathematics, statistics, and computer Science, University of Tehran, Iran)<br>R.Eizadi Alikelaye<br>(Islamic Azad University of QAZVIN, Iran)<br>E-mail: jolany@ut.ac.ir, re.eizadi@gmail.com


#### Abstract

The Arakawa-Kaneko zeta function has been introduced ten years ago by T. Arakawa and M. Kaneko in [22]. In [22], Arakawa and Kaneko have expressed the special values of this function at negative integers with the help of generalized Bernoulli numbers $B^{(k)}$ called poly-Bernoulli numbers. Kim-Kim [4] introduced Multi poly- Bernoulli numbers and proved that special values of certain zeta functions at non-positive integers can be described in terms of these numbers. The study of Multi poly-Bernoulli and Euler numbers and their combinatorial relations has received much attention [2,4,6,7,12,13,14,19,22,27]. In this paper we introduce the generalization of Multi poly-Bernoulli and Euler numbers and consider some combinatorial relationships of the Generalized Multi poly-Bernoulli and Euler numbers of higher order. The present paper deals with Generalization of Multi poly-Bernouli numbers and polynomials of higher order. In 2002, Q. M. Luo and et al (see [11, 23, 24]) defined the generalization of Bernoulli polynomials and Euler numbers. Some earlier results of Luo in terms of generalized Multi poly-Bernoulli and Euler numbers, can be deduced. Also we investigate some relationships between Multi poly-Bernoulli and Euler polynomials.


Key Words: Generalized Multi poly-Bernoulli polynomials, generalized Multi poly-Euler polynomials, stirling numbers, polylogarithm, Multi- polylogarithm.

AMS(2010): 05A10, 05A19

## §1. Introduction

Bernoulli numbers are the signs of a very strong bond between elementary number theory, complex analytic number theory, homotopy theory(the J-homomorphism, and stable homotopy groups of spheres), differential topology(differential structures on spheres), the theory of modular forms(Eisenstein series) and p-adic analytic number theory(the p-adic L-function) of

[^11]Mathematics. For $n \in Z, n \geq 0$, Bernulli numbers $B_{n}$ originally arise in the study of finite sums of a given power of consecutive integers. They are given by $B_{0}=1, B_{1}=-1 / 2, B_{2}=1 / 6, B_{3}=$ $0, B_{4}=-1 / 30, \ldots$, with $B_{2 n+1}=0$ for $n>1$, and

$$
\begin{equation*}
B_{n}=-\frac{1}{n+1} \sum_{m=0}^{n-1}\binom{n+1}{m} B_{m}, \quad n \geq 1 \tag{1}
\end{equation*}
$$

The modern definition of Bernoulli numbers $B_{n}$ can be defined by the contour integral

$$
\begin{equation*}
B_{n}=\frac{n!}{2 \pi i} \oint \frac{z}{e^{z}-1} \frac{d z}{z^{n+1}} \tag{2}
\end{equation*}
$$

where the contour encloses the origin, has radius less than $2 \pi$.
Also Bernoulli polynomials $B_{n}(x)$ are usualy defined(see[1], [4], [5]) by the generating function

$$
\begin{equation*}
G(x, t)=\frac{t e^{x t}}{e^{t}-1}=\sum_{n=0}^{\infty} B_{n}(x) \frac{t^{n}}{n!}, \quad|t|<2 \pi \tag{3}
\end{equation*}
$$

and consequently, Bernoulli numbers $B_{n}(0):=B_{n}$ can be obtained by the generating function

$$
\frac{t}{e^{t}-1}=\sum_{n=0}^{\infty} B_{n} \frac{t^{n}}{n!}
$$

Bernoulli polynomials, first studied by Euler (see[1]), are employed in the integral representation of differentiable periodic functions, and play an important role in the approximation of such functions by means of polynomials (see[14]-[18]).

Euler polynomials $E_{n}(x)$ are defined by the generating function

$$
\begin{equation*}
\frac{2 e^{x t}}{e^{t}+1}=\sum_{n=0}^{\infty} E_{n}(x) \frac{t^{n}}{n!}, \quad|t|<\pi \tag{4}
\end{equation*}
$$

Euler numbers $E_{n}$ can be obtained by the generating function

$$
\begin{equation*}
\frac{2}{e^{t}+1}=\sum_{n=0}^{\infty} E_{n} \frac{t^{n}}{n!} \tag{5}
\end{equation*}
$$

The first four such polynomials, are

$$
\begin{aligned}
& B_{0}(x)=1, B_{1}(x)=x-1 / 2, B_{2}(x)=x^{2}-x+1 / 6 \\
& B_{3}(x)=x^{3}-3 / 2 x^{2}+1 / 2 x, \ldots
\end{aligned}
$$

and

$$
\begin{aligned}
& E_{0}(x)=1, E_{1}(x)=x-1 / 2, E_{2}(x)=x^{2}-x \\
& E_{3}(x)=x^{3}-3 / 2 x^{2}+1 / 4, \ldots
\end{aligned}
$$

Euler polynomials are strictly connected with Bernoulli ones, and are used in the Taylor expansion in a neighborhood of the origin of trigonometric and hyperbolic secant functions.

In the sequel, we list some properties of Bernoulli and Euler numbers and polynomials as well as recurrence relations and identities.

$$
\begin{align*}
B_{n}(x)= & \sum_{k=0}^{n}\binom{n}{k} B_{k} x^{n-k}  \tag{6}\\
E_{n}(x)= & \frac{1}{n+1} \sum_{k=1}^{n+1}\left(2-2^{k+1}\right)\binom{n+1}{k} B_{k} x^{n+1-k}  \tag{7}\\
& B_{n}(x+1)-B_{n}(x)=n x^{n-1}  \tag{8}\\
& E_{n}(x+1)+E_{n}(x)=2 x^{n} \tag{9}
\end{align*}
$$

Lemma 1.1(see[20],[21]) For any integer $n \geq 0$, we have

$$
\begin{align*}
& B_{n}(x+1)=\sum_{k=0}^{n}\binom{n}{k} B_{k}(x)  \tag{10}\\
& E_{n}(x+1)=\sum_{k=0}^{n}\binom{n}{k} E_{k}(x) \tag{11}
\end{align*}
$$

Consequently, from (8), (9) and lemma 1.1, we obtain,

$$
\begin{align*}
& \sum_{k=0}^{n}\binom{n+1}{k} B_{k}(x)=(n+1) x^{n}  \tag{12}\\
& \sum_{k=0}^{n}\binom{n}{k} E_{k}(x)+E_{n}(x)=2 x^{n} \tag{13}
\end{align*}
$$

Lemma 1.3 For any positive integer $n \geq 0$, we have

$$
\begin{align*}
& B_{n}(p x)=p^{n-1} \sum_{r=0}^{p-1} B_{n}\left(x+\frac{r}{p}\right)(p \text { is a positive integer })  \tag{14}\\
& E_{n}(p x)=p^{n} \sum_{r=0}^{p-1}(-1)^{r} E_{n}\left(x+\frac{r}{p}\right)(p \text { is an odd integer }) \tag{15}
\end{align*}
$$

Let us briefly recall $k-t h$ polylogarithm. The polylogarithm is a special function $L i_{k}(z)$, that is defined by the sum

$$
\begin{equation*}
L i_{k}(z):=\sum_{s=1}^{\infty} \frac{z^{s}}{s^{k}} \tag{16}
\end{equation*}
$$

For formal power series $L i_{k}(z)$ is the $k-t h$ polylogarithm if $k \geq 1$, and a rational function if $k \leq 0$. The name of the function come from the fact that it may alternatively be defined as the repeated integral of itself, namely that

$$
\begin{equation*}
L i_{k+1}(z)=\int_{0}^{z} \frac{L i_{k}(t)}{t} d t \tag{17}
\end{equation*}
$$

for integer values of $k$, we have the following explicit expressions

$$
\begin{aligned}
& L i_{1}(z)=-\log (1-z), L i_{0}(z)=\frac{z}{1-z} L i_{-1}(z)=\frac{z}{(1-z)^{2}} \\
& L i_{-2}(z)=\frac{z(1+z)}{(1-z)^{3}}, L i_{-3}(z)=\frac{z\left(1+4 z+z^{2}\right)}{(1-z)^{4}}, \ldots
\end{aligned}
$$

The integral of the Bose-Einstein distribution is expressed in terms of a polylogarithm,

$$
\begin{equation*}
L i_{k+1}(z)=\frac{1}{\Gamma(k+1)} \int_{0}^{\infty} \frac{t^{k}}{\frac{e^{t}}{z}-1} d t \tag{18}
\end{equation*}
$$

Lemma 1.3(see[18]) For $n \in N \cup\{0\}$, we have an explicit formula for $L i_{-n}(z)$ as follow

$$
\begin{align*}
L i_{-n}(z) & =\sum_{k=1}^{n+1} \frac{(-1)^{n+k+1}(k-1)!S(n+1, k)}{(1-z)^{k}}  \tag{19}\\
& (n=1,2, \ldots)
\end{align*}
$$

where $s(n, k)$ are Stirling numbers of the second kind.
Now, we introduce the generalization of $L i_{k}(z)$. Let $r$ be an integer with a value greater than one.

Definition 1.1 Let $k_{1}, k_{2}, \ldots k_{r}$ be integers. The generalization of polylogarithm are defined by

$$
\begin{equation*}
L i_{k_{1}, k_{2}, \ldots, k_{r}}(z)=\sum_{\substack{m_{1}, m_{2}, \ldots, m_{r} \in \mathcal{Z} \\ 0<m_{1}<m_{2}<\ldots<m_{r}}} \frac{z^{m_{r}}}{m_{1}^{k_{1}} \ldots m_{r}^{k_{r}}} \tag{20}
\end{equation*}
$$

The rational numbers $B_{n}^{(k)},(n=0,1,2, \ldots)$ are said to be poly-Bernoulli numbers if they satisfy

$$
\begin{equation*}
\frac{L i_{k}\left(1-e^{-x}\right)}{1-e^{-x}}=\sum_{n=0}^{\infty} B_{n}^{(k)} \frac{x^{n}}{n!} \tag{21}
\end{equation*}
$$

In addition, for any $n \geq 0, B_{n}^{(1)}$ is the classical Bernoulli number, $B_{n}($ see[7], [12]). Also, the rational numbers $H_{n}^{(k)}(u),(n=0,1,2, \ldots)$ are said to be poly-Euler numbers if they satisfy

$$
\begin{equation*}
\frac{L i_{k}\left(1-e^{(1-u)}\right)}{u-e^{t}}=\sum_{n=0}^{\infty} H_{n}^{(k)}(u) \frac{t^{n}}{n!} \tag{22}
\end{equation*}
$$

where $u$ is an algebraic real number and $k \geq 1$.(see[13],[19])
Let us now introduce a generalization of poly-Bernoulli numbers, making use of $L i_{k_{1}, \ldots, k_{r}}(z)$.
Definition 1.1(see[7]) Multi poly-Bernoulli numbers $B_{n}^{\left(k_{0}, \ldots, k_{r}\right)},(n=0,1,2, \ldots)$ are defined for each integer $k_{1}, k_{2}, \ldots, k_{r}$ by the generating series

$$
\begin{equation*}
\frac{L i_{\left(k_{1}, k_{2}, \ldots, k_{r}\right)}\left(1-e^{-t}\right)}{\left(1-e^{-t}\right)^{r}}=\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)} \frac{t^{n}}{n!} \tag{23}
\end{equation*}
$$

By Definition 1.2, the left hand side of (23) is

$$
\begin{equation*}
\frac{1}{1^{k_{1}} 2^{k_{2}} \ldots r^{k_{r}}}+\sum_{\substack{0<m_{1}<\ldots<m_{r} \\ m_{r} \neq r}} \frac{\left(1-e^{-t}\right)^{m_{r}-r}}{m_{1}^{k_{1}} \ldots m_{r}^{k_{r}}} \tag{24}
\end{equation*}
$$

hence we have

$$
\begin{align*}
B_{0}^{\left(k_{1}, . ., k_{r}\right)} & =\frac{1}{1^{k_{1}} 2^{k_{2}} \ldots r^{k_{r}}}  \tag{25}\\
B_{1}^{\left(k_{1}, \ldots k_{r}\right)} & =\sum_{0<m_{1}<\ldots<m_{r}} \frac{1}{m_{1}^{k_{1}} \ldots m_{r-1}^{k_{r-1}}(r+1)^{k_{r}}} \tag{26}
\end{align*}
$$

Definition 1.3 Multi poly-Euler numbers $H_{n}^{\left(k_{1}, \ldots, k_{r}\right)},(n=0,1, \ldots)$ are defined for each integer $k_{1}, \ldots, k_{r}$ by the generating series

$$
\begin{equation*}
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{(1-u)}\right)}{\left(u-e^{t}\right)^{r}}=\sum_{n=0}^{\infty} H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u) \frac{t^{n}}{n!} \tag{27}
\end{equation*}
$$

Kaneko [6] presented the following recurrence formulae for poly-Bernoulli numbers which we state hear.

Theorem 1.1(Kaneko)([2,6,14,22]) For any $k \in Z$ and $n \geq 0$, we have

$$
\begin{align*}
& B_{n}^{(k)}=\frac{1}{n+1}\left\{B_{n}^{(k-1)}-\sum_{m=1}^{n-1}\binom{n}{m-1} B_{m}^{(k)}\right\}  \tag{28}\\
& B_{n}^{(k)}=(-1)^{n} \sum_{k=1}^{n+1} \frac{(-1)^{m-1}(m-1)!\left\{\begin{array}{c}
n \\
m-1
\end{array}\right\}}{m^{k}}  \tag{29}\\
& B_{n}^{(-k)}=\sum_{j=0}^{\min (n, k)}(j!)^{2}\left\{\begin{array}{l}
n+1 \\
j+1
\end{array}\right\}\left\{\begin{array}{l}
k+1 \\
j+1
\end{array}\right\} \quad(n, k \geq 0)  \tag{30}\\
& B_{n}^{(-k)}=B_{k}^{(-n)}(n, k \geq 0)  \tag{31}\\
& B_{n}^{(k)}=\sum_{m=0}^{n}(-1)^{m}\binom{n}{m} B_{n-m}^{(k-1)}\left\{\sum_{l=0}^{m} \frac{(-1)^{l}}{n-l+1}\binom{m}{l} B_{l}^{(1)}\right\} \tag{32}
\end{align*}
$$

where

$$
\left\{\begin{array}{l}
n  \tag{33}\\
m
\end{array}\right\}=\frac{(-1)^{m}}{m!} \sum_{l=0}^{m}(-1)^{l}\binom{m}{l} l^{n} \quad n, m \geq 0
$$

called the second type stirling numbers.
Y.Hamahata and H.Masubuchi in [12], presented the following recurrence formulae for Multi poly-Bernoulli numbers.

Theorem 1.2(H.Masubuchi \& Y.Hamahata) For $n \geq 0$ and $\left(k_{1}, \ldots, k_{r} \in Z\right)$ we have

$$
\begin{align*}
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}=  \tag{34}\\
& (-1)^{n} \sum_{m_{r}=r}^{n+r}\left(\sum_{0<m_{1}<\ldots<m_{r}} \frac{(-1)^{m_{r}-r}\left(m_{r}-r\right)!\left\{\begin{array}{c}
n \\
m_{r}-r
\end{array}\right\}}{m_{1}^{k_{1}} \ldots m_{r}^{k_{r}}}\right)
\end{align*}
$$

If $k_{r} \neq 1$ and $n \geq 1$, then

$$
\begin{equation*}
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}=\frac{1}{n+r}\left\{B_{n}^{\left(k_{1}, \ldots, k_{r-1}, k_{r}-1\right)}-\sum_{m=1}^{n-1}\binom{n}{m-1} B_{m}^{\left(k_{1}, \ldots, k_{r}\right)}\right\} \tag{35}
\end{equation*}
$$

If $k_{r}=1$ and $n \geq 1$,then

$$
\begin{align*}
& B_{n}^{\left(k_{1}, \ldots, k_{r-1}, 1\right)}=\frac{1}{n+r}  \tag{36}\\
& \left\{B_{n}^{\left(k_{1}, \ldots, k_{r}-1\right)}-\sum_{m=0}^{n-1}(-1)^{n-m}\left\{r\binom{n}{m}+\binom{n}{m-1}\right\} B_{m}^{\left(k_{1}, \ldots, k_{r-1}, 1\right)}\right\}
\end{align*}
$$

Also, they proved (see[1]) if

$$
\begin{equation*}
B[r]_{n}^{(k)}=B_{n}^{(\overbrace{0, \ldots, 0}^{r-1}}, k) \tag{37}
\end{equation*}
$$

then for $n, k \geq 0$, we have

$$
\begin{equation*}
B[r]_{n}^{(-k)}=B[r]_{k}^{(-n)} \tag{38}
\end{equation*}
$$

In [23], [24], Q.M.Luo, F.Oi and L.Debnath defined the generalization of Bernoulli and Euler polynomials $B_{n}(x, a, b, c)$ and $E_{n}(x, a, b, c)$ respectively, which are expressed as follows

$$
\begin{align*}
\frac{t}{b^{t}-a^{t}} c^{x t} & =\sum_{k=0}^{\infty} B_{k}(x, a, b, c) \frac{t^{k}}{k!}  \tag{39}\\
\frac{2 c^{x t}}{b^{t}+a^{t}} & =\sum_{k=0}^{\infty} E_{k}(x, a, b, c) \frac{t^{k}}{k!} \tag{40}
\end{align*}
$$

In this paper, by the method of Q.M.Luo and et al [11], we give some properties on generalized Multi poly-Bernoulli and Euler polynomials

Definition 1.4 Let $a, b>0$ and $a \neq b$. The generalized Multi poly-Bernoulli numbers $B_{n}^{\left(k_{1}, \cdots, k_{r}\right)}(a, b)$, the generalized Multi poly-Bernoulli polynomials

$$
B_{n}^{\left(k_{1}, \cdots, k_{r}\right)}(x, a, b) \quad \text { and } \quad B_{n}^{\left(k_{1}, \cdots, k_{r}\right)}(x, a, b, c)
$$

are defined by the following generating functions, respectively;

$$
\begin{align*}
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-t}\right)}{\left(b^{t}-a^{-t}\right)^{r}} & =\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b) \frac{t^{n}}{n!}, \quad|t|<\frac{2 \pi}{|\ln a+\ln b|}  \tag{41}\\
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-t}\right)}{\left(b^{t}-a^{-t}\right)^{r}} e^{r x t} & =\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b) \frac{t^{n}}{n!}, \quad|t|<\frac{2 \pi}{|\ln a+\ln b|}  \tag{42}\\
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-t}\right)}{\left(b^{t}-a^{-t}\right)^{r}} c^{r x t} & =\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c) \frac{t^{n}}{n!}, \quad|t|<\frac{2 \pi}{|\ln a+\ln b|} \tag{43}
\end{align*}
$$

Definition 1.5 Let $a, b>0$, and $a \neq b$, the generalized Multi poly-Euler numbers $H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; a, b)$, the generalized multi poly-Euler polynomial $H_{n}^{k_{1}, \ldots, k_{r}}(x ; u, a, b)$ and $H_{n}^{k_{1}, \ldots, k_{r}}(x ; u, a, b, c)$ are defined by the following generating functions, respectively,

$$
\begin{align*}
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{(1-u)}\right)}{\left(u a^{-t}-b^{t}\right)^{r}} & =\sum_{n=0}^{\infty} H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u, a, b) \frac{t^{n}}{n!}, \quad|t|<\frac{2 \pi}{|\ln a+\ln b|}  \tag{44}\\
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{(1-u)}\right)}{\left(u a^{-t}-b^{t}\right)^{r}} e^{r x t} & =\sum_{n=0}^{\infty} H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; u, a, b) \frac{t^{n}}{n!}, \quad|t|<\frac{2 \pi}{|\ln a+\ln b|}  \tag{45}\\
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{(1-u)}\right)}{\left(u a^{-t}-b^{t}\right)^{r}} c^{r x t} & =\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; u, a, b, c) \frac{t^{n}}{n!}, \quad|t|<\frac{2 \pi}{|\ln a+\ln b|} \tag{46}
\end{align*}
$$

## §2. Main Theorems

In this section, we introduce our main results. We give some theorems and corollaries which are related to generalized Multi poly-Bernoulli numbers and generalized Multi poly-Euler polynomials. We present some recurrence formulae for generalized Multi-poly-Bernoulli and Euler polynomials.

Theorem 2.1 Let $a, b>0$ and $a \neq b$, we have

$$
\begin{equation*}
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b)=B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(\frac{-\ln b}{\ln a+\ln b}\right)(\ln a+\ln b)^{n} \tag{47}
\end{equation*}
$$

proof By applying Definition 1.4, we have

$$
\begin{aligned}
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-x}\right)}{\left(b^{x}-a^{-x}\right)^{r}} & =\sum_{n=0}^{\infty} \frac{B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b)}{n!} x^{n} \\
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-x}\right)}{\left(b^{x}-a^{-x}\right)^{r}} & =\frac{1}{b^{x r}}\left(\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{-x \ln a b}\right)}{\left(1-e^{-x \ln a b}\right)^{r}}\right) \\
& =e^{-x r \ln b}\left(\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{-x \ln a b}\right)}{\left(1-e^{-x \ln a b}\right)^{r}}\right)
\end{aligned}
$$

So, we get

$$
\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{-x \ln a b}\right)}{\left(b^{x}-a^{-x}\right)^{r}}=\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(\frac{-\ln b}{\ln a+\ln b}\right)(\ln a+\ln b)^{n} \frac{x^{n}}{n!}
$$

Therefore, by comparing the coefficients of $\frac{t^{n}}{n!}$ on both sides, proof will be complete

$$
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b)=B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(\frac{-\ln b}{\ln a+\ln b}\right)(\ln a+\ln b)^{n}
$$

The generalized Multi poly-Bernoulli and Euler numbers process a number of interesting properties which we state here

Theorem 2.2 Let $a, b>0$ and $a \neq b$. For real algebraic $u$ we have

$$
\begin{equation*}
H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; a, b)=H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(u ; \frac{\ln a}{\ln a+\ln b}\right)(\ln a+\ln b)^{n} \tag{48}
\end{equation*}
$$

Next, we investigate a strong relationships between $B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b)$ and $B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}$.
Theorem 2.3 Let $a, b>0, a \neq b a n d a>b>0$, we have

$$
\begin{equation*}
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b)=\sum_{i=0}^{j}(-r)^{j-i}(\ln a+\ln b)^{i}(\ln b)^{j-i}\binom{j}{i} B_{i}^{\left(k_{1}, \ldots, k_{r}\right)} . \tag{49}
\end{equation*}
$$

By applying Definition 1.4, we have

$$
\begin{aligned}
& \frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-x}\right)}{\left(b^{x}-a^{-x}\right)^{r}}=\frac{1}{b^{x r}} \frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-x}\right)}{\left(1-e^{-x \ln a b}\right)^{r}} \\
& =\left(\sum_{k=0}^{\infty} \frac{(\ln b)^{k}}{k!} x^{k} r^{k}(-1)^{k}\right)\left(\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(\ln a+\ln b)^{n} \frac{x^{n}}{n!}\right) \\
& =\sum_{j=0}^{\infty}\left(\sum_{i=0}^{j}(-r)^{j-i} \frac{B_{j}^{\left(k_{1}, \ldots, k_{r}\right)}(\ln a+\ln b)^{i}(\ln b)^{j-i}}{i!(j-i)!} x^{j}\right)
\end{aligned}
$$

By comparing the coefficient of $\frac{t^{n}}{n!}$ on both sides, we get.

$$
B_{j}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b)=\sum_{i=0}^{j}(-r)^{j-i}(\ln a+\ln b)^{i}(\ln b)^{j-i}\binom{j}{i} B_{i}^{\left(k_{1}, \ldots, k_{r}\right)}
$$

By the same method proceeded in the proof of Theorem 2.3, we obtained similar relations for $H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; a, b)$ and $H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}$.

Theorem 2.4 Let $a, b>0$, and $b>a>0$. For algebraic real number $u$, we have

$$
\begin{equation*}
H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; a, b)=\sum_{i=0}^{n} r^{i}(\ln a+\ln b)^{i}(\ln a)^{n-i}\binom{n}{i} H_{i}^{\left(k_{1}, \ldots, k_{r}\right)} \tag{50}
\end{equation*}
$$

Theorem 2.5 Let $x \in R$ and conditions of Theorem 2.3 holds true, then we get

$$
\begin{array}{r}
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c)=\sum_{l=0}^{n}\binom{n}{l} r^{n-l}(\ln c)^{n-l} B_{l}^{\left(k_{1}, \ldots, k_{r}\right)}(a, b) x^{n-l} \\
H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; x, a, b, c)=\sum_{l=0}^{n}\binom{n}{l} r^{n-l}(\ln c)^{n-l} H_{l}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; a, b) x^{n-l} \tag{52}
\end{array}
$$

Proof By applying Definitions 1.4 and 1.5 , proof will be complete.

Theorem 2.6 Let conditions of Theorem 2.5 holds true, we obtain

$$
\begin{align*}
& H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; x, a, b, c)= \\
& \sum_{k=0}^{n}\binom{n}{k} r^{n-k}(\ln c)^{n-k} H_{k}^{\left(k_{1}, \ldots, k_{r}\right)}\left(u, \frac{\ln a}{\ln a+\ln b}\right)(\ln a+\ln b)^{k} x^{n-k}  \tag{53}\\
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c)= \\
& \sum_{k=0}^{n}\binom{n}{k} r^{n-k}(\ln c)^{n-k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}\left(\frac{-\ln b}{\ln a+\ln b}\right)(\ln a+\ln b)^{k} x^{n-k} \tag{54}
\end{align*}
$$

Proof By applying Theorems 2.1 and 2.5, we get (53), and Obviously, the result of (54) is similar with (53).

Theorem 2.7 Let conditions of Theorem 2.5 holds true, then we get

$$
\begin{align*}
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c)=  \tag{55}\\
& =\sum_{k=0}^{n} \sum_{j=0}^{k}(-1)^{k-j}\binom{n}{k}\binom{k}{j} r^{n-k}(\ln c)^{n-k}(\ln b)^{k-j}(\ln a+\ln b)^{j} B_{j}^{\left(k_{1}, \ldots, k_{r}\right)} x^{n-k} \\
& H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c)=  \tag{56}\\
& =\sum_{k=0}^{n} \sum_{j=0}^{k}\binom{n}{k}\binom{k}{j} r^{n-k}(\ln c)^{n-k}(\ln a)^{k-j}(\ln c+\ln b)^{j} H_{j}^{\left(k_{1}, \ldots, k_{r}\right)} x^{n-k} \\
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x+1 ; a, b, c)=B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(x ; a c, \frac{b}{c}, c\right)  \tag{57}\\
& H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u, 1-x, a c, b, c)=B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(u,-x, a c, \frac{b}{c}, c\right)  \tag{58}\\
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x+y ; a, b, c)=\sum_{k=0}^{n}\binom{n}{k} r^{n-k}(\ln c)^{n-k} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c) y^{n-k}  \tag{59}\\
& H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; x+y, a, b, c)=\sum_{k=0}^{n}\binom{n}{k} r^{n-k}(\ln c)^{n-k} H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c) y^{n-k} \tag{60}
\end{align*}
$$

Proof We only prove (59) and (55)-(60) can be derived by Definitions 1.4 and 1.5.

$$
\begin{aligned}
& \frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-t}\right)}{\left(b^{t}-a^{-t}\right)^{r}} c^{(x+y) r t}=\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x+y, a, b, c) \frac{t^{n}}{n!} \\
& =\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-t}\right)}{\left(b^{t}-a^{-t}\right)^{r}} c^{x r t} . c^{y r t} \\
& =\left(\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c) \frac{t^{n}}{n!}\right)\left(\sum_{i=0}^{n} \frac{y^{i}(\ln c)^{i} r^{i}}{i!} t^{i}\right) \\
& =\sum_{n=0}^{\infty}\left(\sum_{k=0}^{n}\binom{n}{k} r^{n-k} y^{n-k}(\ln c)^{n-k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c)\right) \frac{t^{n}}{n!}
\end{aligned}
$$

So by comparing the coefficients of $\frac{t^{n}}{n!}$ in the two expressions, we obtain the desired result 2.13 .

Theorem 2.8 By the same method proceeded in the proof of previous Theorems, we find similar relations for $B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(t)$ and $H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u, t)$.

$$
\begin{align*}
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(t)=B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(e^{1+t}, e^{-t}\right)  \tag{61}\\
& H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u, t)=H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(u ; e^{t}, e^{1-t}\right) \tag{62}
\end{align*}
$$

Now, we present formulae which show a deeper motivation of generalized poly-Bernoulli and Euler polynomials.

Theorem 2.9 Let $x, y \in R$ and conditions of Theorem 2.5 holds true, we get

$$
\begin{align*}
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x, a, b, c) & =(\ln a+\ln b)^{n} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(\frac{-\ln b+x \ln c}{\ln a+\ln b}\right)  \tag{63}\\
H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(u ; x, a, b, c) & =H_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(u ; \frac{\ln a+x \ln c}{\ln a+\ln b}\right) \tag{64}
\end{align*}
$$

Proof We can write

$$
\begin{aligned}
\sum_{n=0}^{\infty} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c) \frac{t^{n}}{n!} & =\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-t}\right)}{\left(b^{t}-a^{-t}\right)^{r}} c^{x r t} \\
& =\frac{1}{b^{r t}} \frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-(a b)^{-t}\right)}{\left(1-(a b)^{-t}\right)^{r}} c^{x r t} \\
& =e^{r(-\ln b+x \ln c) t}\left(\frac{L i_{\left(k_{1}, \ldots, k_{r}\right)}\left(1-e^{-t \ln a b}\right)}{\left(1-e^{-t \ln a b}\right)^{r}}\right)
\end{aligned}
$$

By comparing the coefficients of $\frac{t^{n}}{n!}$ on both sides, we get

$$
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x ; a, b, c)=(\ln a+\ln b)^{n} B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}\left(\frac{-\ln b+x \ln c}{\ln a+\ln b}\right)
$$

GI-Sang Cheon and H.M.Srivastava in [8],[10] investigated the classical relationship between Bernoulli and Euler polynomials . Now we present a relationship between generalized Multi poly-Bernoulli and generalized Euler polynomials. The following relation (65) are given by Q.M.Luo, So by applying this recurrence formula, we obtain Theorem 2.10,

$$
\begin{equation*}
E_{k}(x+1,1, b, b)+E_{k}(x, 1, b, b)=2 x^{k}(\ln b)^{k} \tag{65}
\end{equation*}
$$

Theorem 2.10 Let $a, b>0$, we have

$$
\begin{gather*}
B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x+y ; a, b)=  \tag{66}\\
\frac{1}{2} \sum_{k=0}^{n}\binom{n}{k}\left[B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y, a, b)+B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y+1, a, b)\right] r^{n-k} E_{n-k}(x, 1, b, b)
\end{gather*}
$$

Proof We know

$$
\begin{aligned}
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x+y ; 1, b, b)=\sum_{k=0}^{n}\binom{n}{k} r^{n-k}(\ln b)^{n-k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y ; 1, b, b) x^{n-k} \\
& E_{k}(x+y, 1, b, b)+E_{k}(x, 1,, b, b)=2 x^{k}(\ln b)^{k}
\end{aligned}
$$

So, we obtain

$$
\begin{aligned}
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x+y, 1, b, b)=\frac{1}{2} \sum_{k=0}^{n}\binom{n}{k} r^{n-k}(\ln b)^{n-k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y ; 1, b, b) \times \\
& {\left[\frac{1}{(\ln b)^{n-k}}\left(E_{n-k}(x ; 1, b, b)+E_{n-k}(x+1,1, b, b)\right)\right]} \\
& =\frac{1}{2} \sum_{k=0}^{n}\binom{n}{k} r^{n-k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y ; 1, b, b) \times \\
& {\left[\left(E_{n-k}(x ; 1, b, b)+\sum_{j=0}^{n-k}\binom{n-k}{j} E_{j}(x, 1, b, b)\right)\right]} \\
& \quad=\frac{1}{2} \sum_{k=0}^{n}\binom{n}{k} r^{n-k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y ; 1, b, b) E_{n-k}(x ; 1, b, b) \\
& \quad+\frac{1}{2} \sum_{j=0}^{n}\binom{n}{j} r^{n-k} E_{j}(x ; 1, b, b) \sum_{k=0}^{n-j}\binom{n-j}{k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y ; 1, b, b) \\
& \quad=\frac{1}{2} \sum_{k=0}^{n}\binom{n}{k} r^{n-k} B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y ; 1, b, b) E_{n-k}(x ; 1, b, b) \\
& \quad+\frac{1}{2} \sum_{j=0}^{n}\binom{n}{j} r^{n-k} B_{n-j}^{\left(k_{1}, \ldots, k_{r}\right)}(y+1 ; 1, b, b) E_{j}(x ; 1, b, b)
\end{aligned}
$$

So we have

$$
\begin{aligned}
& B_{n}^{\left(k_{1}, \ldots, k_{r}\right)}(x+y ; a, b) \\
& =\frac{1}{2} \sum_{k=0}^{n}\binom{n}{k}\left[B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y, a, b)+B_{k}^{\left(k_{1}, \ldots, k_{r}\right)}(y+1, a, b)\right] r^{n-k} E_{n-k}(x, 1, b, b)
\end{aligned}
$$

Therefore we obtain the desired result (66).
The following corollary is a straightforward consequence of Theorem 2.10.
Corollary 2.1(see [8],[10]) In Theorem 2.10, if we set $r=1, k=1$ and $b=e$, we obtain

$$
\begin{equation*}
B_{n}(x)=\sum_{\substack{k=0 \\ k \neq 1}}^{n}\binom{n}{k} B_{k} E_{n-k}(x) . \tag{67}
\end{equation*}
$$

Further work: In [25], Jang et al. gave new formulae on Genocchi numbers. They defined poly-Genocchi numbers to give the relation between Genocchi numbers, Euler numbers, and poly-Genocchi numbers. After Y. Simsek [26], gave a new generating functions which produce Genocchi zeta functions. So by applying a similar method of Kim-Kim [4], we can introduce generalized Genocchi Zeta functions and next define Multi poly-Genocchi numbers and obtain several properties in this area.

Acknowledgments: The authors wishes to express his sincere gratitude to the referee for his/her valuable suggestions and comments and Professor Mohammad Maleki for his cooperations and helps.

## References

[1] M.Abramowits, I. A.Stegun, Handbook of Mathematical functions with formulas, Graphs and Mathematical tables,National Bureau of standards, Washington, DC, 1964.
[2] T.Arakawa and M.Kaneko, On poly-Bernoulli numbers, Comment Math. univ. st. Pauli, 48, (1999), 159-167.
[3] B. N.Gue and F.Qi, Generalization of Bernoulli polynomials, I. J. Math. Ed. Sci. Tech, 33, (2002), N0 3, 428-431
[4] M.-S.-Kim and T.Kim, An explicit formula on the generalized Bernoulli number with other n, Indian. J. Pure and applied Math., 31, (2000), 1455-1466.
[5] Hassan jolany and M.R.Darafsheh, Some other remarks on the generalization of Bernoulli and Euler numbers, Journal's name??, Vol. 5 (2009), No. 3, 118-129.
[6] M.Kaneko, Poly-bernoulli numbers, Journal de Theorides de bordeaux, 9(1997), 221-228
[7] Y.Hamahata, H.Masubuch, Special multi-poly-Bernoulli numbers, Journal of Integer sequences, Vol 10, (2007).
[8] H. M.Srivastava and A.Pinter, Remarks on some relationships between the Bernoulli and Euler polynomials, Applied. Math. Letter, 17,(2004)375-380.
[9] Chad Berwbaker, A combinatorial interpretion of the poly-Bernoulli numbers and two fermat analogues, Integers Journal, 8, (2008).
[10] GI-Sang Cheon, A note on the Bernoulli and Euler polynomials, Applied. Math. Letter, 16, (2003), 365-368.
[11] Q. M.Lue, F.Oi and L.Debnath, Generalization of Euler numbers and polynomials, Int. J. Math. Sci., (2003) 3893-3901.
[12] Y.Hamahata, H.Masubuchi, Recurrence formulae for multi-poly-Bernoulli numbers, Integers Journal, 7(2007).
[13] Jin-Woo Son and Min Soo kim, On Poly-Eulerian numbers, Bull.Korean Math. Sco., 1999 36:47-61.
[14] M.Kaneko, N.Kurokawa, M.Wakayama, A variation of Euler's approach to the Riemann zeta function, Kyushu J.Math., 57(2003) 175-192.
[15] M.Ward, A calculus of sequences, Amer.J.Math., 58(1936) 255-266.
[16] H.Tsumura, A note on q-analogues of the Dirichlet series and q-Bernoulli numbers, J.Number Theory, 39(1991) 251-256.
[17] T.Kim, On a q-analogue of the p-adic log gamma functions and related integrals, J.Number Theory, 76(1999)320-329.
[18] H.M.Srivastava and J.Choi, Series Associated with The zeta and Related functions, Kluwer Academic Publisher, Dordrechet,Boston, London, 2001.
[19] Hassan Jolany and M.R.Darafsheh, Generalization on poly-Bernoulli numbers and polynomials, International J.Math. Combin., Vol. 2 (2010), 07-14.
[20] Q-M-Luo, H.M.Srivastava, Some relationships between the Apostol-Bernoulli and ApostolEuler polynomials, Comput.Math.Appl., 51(2006) 631-642.
[21] L.Carlitz,q-Bernoulli and Eulerian numbers, Trans.Amer.math.Soc., 76 (1954) 332-350
[22] T.Arakawa, M.Kaneko, Multiple zeta values, Poly-Bernoulli numbers and related zeta functions, Nagoya Math. J., 153 (1999), 189-209.
[23] Q.M.Luo, B. N.Guo, F.Qi, and L.Debnath,Generalization of Bernoulli numbers and polynomials, IJMMS, Vol. 2003, Issue 59, 2003, 3769-3776.
[24] Q.M.Luo, F.Qi, and L.Debnath, Generalizations of Euler numbers and polynomials, IJMMS, Vol. 2003, Issue 61, 2003, 3893-3901.
[25] L.C.Jang, T.Kim, D.H. Lee and D.W. Park, An application of polylogarithms in the analogs of Genocchi numbers, Notes on Number Theory and Discrete Mathematics, Vol.7, No. 3, pp. 65-69, 2001.
[26] Y.Simsek, q-Hardy-Berndt type sums associated with q-Genocchi type zeta and l-functions, http://arxiv.org/abs/0710.5681v1.
[27] Hassan Jolany and Hamed Faramarzi, Generalizations on poly-Eulerian numbers and polynomials, Scientia Magna, Vol. 6 (2010), No. 1, 9-18.

# Bounds on the Largest of 

 Minimum Degree Laplician Eigenvalues of a GraphChandrashekar Adiga and C. S. Shivakumar Swamy<br>(Department of Studies in Mathematics, University of Mysore, Manasagangotri, Mysore 570 006, India)<br>E-mail: c_ adiga@hotmail.com, cskswamy@gmail.com


#### Abstract

In this paper we give three upper bounds for the largest of minimum degree Laplacian eigenvalues of a graph and also obtain a lower bound for the same.


Key Words: Minimum degree matrix, minimum degree Laplacian eigenvalues.
AMS(2010): 05C50

## §1. Introduction

Let $G=(V, E)$ be a simple, connected graph with vertex set $V=\left\{v_{1}, v_{2}, \ldots, v_{n}\right\}$ and edge set $E=\left\{e_{1}, e_{2}, \ldots, e_{m}\right\}$. Assume that the vertices are ordered such that $d_{1} \geq d_{2} \geq \ldots \geq d_{n}$, where $d_{i}$ is the degree of $v_{i}$ for $i=1,2, \ldots, n$. The energy of $G$ was first defined by I.Gutman [5] in 1978 as the sum of the absolute values of its eigenvalues. The energy of a graph has close links to Chemistry ( see for instance [6]). The $n \times n$ matrix $m(G)=\left(d_{i j}\right)$ is called the minimum degree matrix of $G$, where

$$
d_{i j}= \begin{cases}\min \left\{d_{i}, d_{j}\right\} & \text { if } v_{i} \text { and } v_{j} \text { are adjacent } \\ 0 & \text { otherwise }\end{cases}
$$

This was introduced and studied in [1]. The characteristic polynomial of the minimum degree matrix $m(G)$ is defined by

$$
\begin{align*}
\phi(G ; \lambda) & =\operatorname{det}(\lambda I-m(G)) \\
& =\lambda^{n}+c_{1} \lambda^{n-1}+c_{2} \lambda^{n-2}+\ldots+c_{n-1} \lambda+c_{n} \tag{1.1}
\end{align*}
$$

where $I$ is the unit matrix of order $n$. The minimum degree Laplacian matrix of $G$ is $L(G)=$ $D(G)-m(G)$, where $D(G)=\operatorname{diag}\left(d_{1}, d_{2}, \ldots, d_{n}\right) . L(G)$ is a real, symmetric matrix. The minimum degree Laplacian eigenvalues $\mu_{1}, \mu_{2}, \ldots, \mu_{n}$ of the graph $G$, assumed in the non increasing order, are the eigenvalues of $L(G)$. The Laplacian matrix of $G$ is $L_{1}(G)=D(G)-A(G)$, where $A(G)$ is the adjacency matrix of $G$. The eigenvalues of the laplacian matrix $L_{1}(G)$ are important in graph theory, because they have relations to numerous graph invariants including connectivity, expanding property, isoperimetric number, independence number, genus, diameter, mean distance, and bandwidth-type parameters of a graph(see, for example, [2,3,9,10]). In

[^12]many applications one needs good bounds for the largest Laplacian eigenvalue (see for instance [ $2,3,9,10]$ ). In this paper, we give three upper bounds and a lower bound for $\mu_{1}$ the largest of minimum degree Laplacian eigenvalues of a graph.

## §2. Main Results

In this section, we will give three upper bounds for $\mu_{1}$ the largest of minimum degree Laplacian eigenvalues of a graph. We employ the following theorem to prove one of our main results.

Theorem 2.1([4]) Let $G$ be a simple graph with $n$ vertices and $m$ edges, and let $\Pi=\left(d_{1}, d_{2}, \ldots, d_{n}\right)$ be the degree sequence of $G$. Then,

$$
d_{1}^{2}+d_{2}^{2}+\ldots+d_{n}^{2} \leq m\left(\frac{2 m}{n-1}+n-2\right)
$$

Theorem 2.2 Let $G$ be a connected graph with $n$ vertices and $m$ edges. Then

$$
\mu_{1} \leq \frac{2 m+\sqrt{(n-1)\left[n\left(2\left|c_{2}\right|+m\left(\frac{2 m}{n-1}+n-2\right)\right)-4 m^{2}\right]}}{n}
$$

where $c_{2}$ is the coefficient of $\lambda^{n-2}$ in $\operatorname{det}(\lambda I-m(G))$.
Proof Clearly

$$
\begin{gather*}
\mu_{1}+\mu_{2}+\ldots+\mu_{n}=\operatorname{Trace}[L(G)]=\sum_{v \in V(G)} d_{v}  \tag{2.1}\\
\mu_{1}^{2}+\mu_{2}^{2}+\ldots+\mu_{n}^{2}=2\left|c_{2}\right|+\sum_{i=1}^{n} d_{i}^{2} \tag{2.2}
\end{gather*}
$$

By Cauchy-Schwarz inequality, we have

$$
\begin{equation*}
\left(\sum_{i=1}^{n} a_{i} b_{i}\right)^{2} \leq\left(\sum_{i=1}^{n} a_{i}^{2}\right)\left(\sum_{i=1}^{n} b_{i}^{2}\right) \tag{2.3}
\end{equation*}
$$

Putting $a_{i}=1$ and $b_{i}=\mu_{i}$ for $i=2, \ldots, n$ in (2.3), we get

$$
\left(\sum_{i=1}^{n} \mu_{i}-\mu_{1}\right)^{2} \leq(n-1)\left(\sum_{i=1}^{n} \mu_{i}^{2}-\mu_{1}^{2}\right)
$$

Using (2.1) and (2.2) in above inequality, we obtain

$$
\left(\sum_{v \in V(G)} d_{v}-\mu_{1}\right)^{2} \leq(n-1)\left[2\left|c_{2}\right|+\sum_{i=1}^{n} d_{i}^{2}\right]-(n-1) \mu_{1}^{2}
$$

After some simplifications, we deduce that

$$
\left(n \mu_{1}-\sum_{v \in V(G)} d_{v}\right)^{2}+(n-1)\left(\sum_{v \in V(G)} d_{v}\right)^{2} \leq n(n-1)\left[2\left|c_{2}\right|+\sum_{i=1}^{n} d_{i}^{2}\right]
$$

$$
\text { i.e., } \quad n \mu_{1}-\sum_{v \in V(G)} d_{v} \leq \sqrt{(n-1)\left[n\left(2\left|c_{2}\right|+\sum_{i=1}^{n} d_{i}^{2}\right)-\left(\sum_{i=1}^{n} d_{i}\right)^{2}\right]}
$$

Therefore

$$
\begin{equation*}
\mu_{1} \leq \frac{\sum_{i=1}^{n} d_{i}+\sqrt{(n-1)\left[n\left(2\left|c_{2}\right|+\sum_{i=1}^{n} d_{i}^{2}\right)-\left(\sum_{i=1}^{n} d_{i}\right)^{2}\right]}}{n} \tag{2.4}
\end{equation*}
$$

Employing Theorem 2.1 and $\sum_{i=1}^{n} d_{i}=2 m$ in (2.4), we see that

$$
\mu_{1} \leq \frac{2 m+\sqrt{(n-1)\left[n\left(2\left|c_{2}\right|+m\left(\frac{2 m}{n-1}+n-2\right)\right)-4 m^{2}\right]}}{n}
$$

This completes the proof.
The following theorem gives another type of upper bound for $\mu_{1}$.
Theorem 2.3 Let $G$ be connected graph with $n$ vertices and $m$ edges. Then

$$
\mu_{1} \leq \sqrt{2 d_{1}^{2}+4 m-2 d_{n}^{3}\left(n-d_{1}\right)}
$$

Proof Suppose that $X=\left(x_{1}, x_{2}, x_{3}, \ldots, x_{n}\right)^{T}$ be an eigenvector with unit length corresponding to $\mu_{1}$. Then

$$
L(G) X=\mu_{1} X
$$

Hence, for $u \in V(G)$,

$$
\mu_{1} x_{u}=d_{u} x_{u}-\sum_{\substack{v \in V(G) \\ v \neq u}} d_{u v} x_{v}
$$

Here $x_{u}$ we mean $x_{i}$ if $u=v_{i}$. Therefore

$$
\begin{equation*}
\mu_{1} x_{u}=\sum_{v u \in E(G)}\left(x_{u}-\min \left(d_{u}, d_{v}\right) x_{v}\right) \tag{2.5}
\end{equation*}
$$

By Cauchy-Schwarz inequality, we have

$$
\begin{aligned}
\mu_{1}^{2} x_{u}^{2} & \leq\left(\sum_{v u \in E(G)} 1^{2}\right)\left(\sum_{v u \in E(G)}\left(x_{u}-\min \left(d_{u}, d_{v}\right) x_{v}\right)^{2}\right) \\
& =d_{u}\left[\sum_{v u \in E(G)} x_{u}^{2}+\sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2}-2 x_{u} \min \left(d_{u}, d_{v}\right) x_{v}\right] .
\end{aligned}
$$

Observe that

$$
\begin{equation*}
-2 x_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right) x_{v} \leq d_{u} x_{u}^{2}+\sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2} \tag{2.6}
\end{equation*}
$$

Hence,

$$
\begin{gather*}
\mu_{1}^{2} x_{u}^{2} \leq d_{u}\left[\sum_{v u \in E(G)} x_{u}^{2}+\sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2}+d_{u} x_{u}^{2}+\sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2}\right] . \\
\text { i.e., }  \tag{2.7}\\
\mu_{1}^{2} x_{u}^{2} \leq 2 d_{u}^{2} x_{u}^{2}+2 d_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2} .
\end{gather*}
$$

Consequently,

$$
\begin{aligned}
\mu_{1}^{2}= & \mu_{1}^{2} \sum_{u \in V(G)} x_{u}^{2} \\
& \leq \sum_{u \in V(G)}\left[2 d_{u}^{2} x_{u}^{2}+2 d_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2}\right] \\
= & 2 \sum_{u \in V(G)} d_{u}^{2} x_{u}^{2}+2 \sum_{u \in V(G)} d_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2} .
\end{aligned}
$$

Thus

$$
\begin{equation*}
\mu_{1}^{2} \leq 2 d_{1}^{2}+2 \sum_{u \in V(G)} d_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2} \tag{2.8}
\end{equation*}
$$

Now let $v \nsim u$ mean that $u$ and $v$ are not adjacent. Then

$$
\begin{aligned}
& \sum_{u \in V(G)} d_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2} \\
& =\sum_{u \in V(G)} d_{u}\left[1-\sum_{v \nsim u} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2}\right]=2 m-\sum_{u \in V(G)} d_{u} \sum_{v \nsim u} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2} \\
& =2 m-\left(\sum_{u \in V(G)} d_{u} \min \left(d_{u}, d_{v}\right)^{2} x_{u}^{2}+\sum_{u \in V(G)} d_{u} \sum_{v \nsim u, v \neq u} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2}\right) \\
& \leq 2 m-\left(d_{n}^{2} \sum_{u \in V(G)} d_{u} x_{u}^{2}+\sum_{u \in V(G)} d_{n} \sum_{v \nsim u, v \neq u} d_{n}^{2} x_{v}^{2}\right) \\
& = \\
& =2 m-\left(d_{n}^{2} \sum_{u \in V(G)} d_{u} x_{u}^{2}+\sum_{u \in V(G)} d_{n}^{3}\left(n-d_{u}-1\right) x_{u}^{2}\right) \\
& = \\
& \leq 2 m-\left(d_{n}^{2} \sum_{u \in V(G)} d_{u} x_{u}^{2}+d_{n}^{3} \sum_{u \in V(G)} n x_{u}^{2}-d_{n}^{3} \sum_{u \in V(G)} d_{u} x_{u}^{2}-d_{n}^{3} \sum_{u \in V(G)} x_{u}^{2}\right) \\
& = \\
& =2 m-d_{n}^{3}\left(n-d_{1}\right) .
\end{aligned}
$$

Hence, employing this in (2.8) we have

$$
\mu_{1}^{2} \leq 2 d_{1}^{2}+4 m-2 d_{n}^{3}\left(n-d_{1}\right) .
$$

Therefore

$$
\mu_{1} \leq \sqrt{2 d_{1}^{2}+4 m-2 d_{n}^{3}\left(n-d_{1}\right)}
$$

Theorem 2.4 Let $G$ be a connected graph then

$$
\mu_{1} \leq \max \left(\sqrt{2\left(d_{u}^{2}+d_{1}^{2} m_{u} d_{u}\right)} \quad: \quad u \in V(G)\right)
$$

Proof From (2.7) we have

$$
\mu_{1}^{2} x_{u}^{2} \leq 2 d_{u}^{2} x_{u}^{2}+2 d_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2}
$$

Thus

$$
\begin{aligned}
\mu_{1}^{2} \sum_{u \in V(G)} x_{u}^{2} & \leq 2 \sum_{u \in V(G)} d_{u}^{2} x_{u}^{2}+2 \sum_{u \in V(G)} d_{u} \sum_{v u \in E(G)} \min \left(d_{u}, d_{v}\right)^{2} x_{v}^{2} \\
& \leq 2 \sum_{u \in V(G)} d_{u}^{2} x_{u}^{2}+2 d_{1}^{2} \sum_{u \in V(G)} d_{u} \sum_{v u \in E(G)} x_{v}^{2} \\
& =2\left[\sum_{u \in V(G)} d_{u}^{2} x_{u}^{2}+d_{1}^{2} \sum_{u \in V(G)} x_{u}^{2} \sum_{v u \in E(G)} d_{v}\right] \\
& =2\left[\sum_{u \in V(G)} d_{u}^{2} x_{u}^{2}+d_{1}^{2} \sum_{u \in V(G)} x_{u}^{2} m_{u} d_{u}\right]
\end{aligned}
$$

where $m_{u}=$ average degree of the vertices adjacent to $u$.
So,

$$
\mu_{1} \leq \sqrt{2 \sum_{u \in V(G)}\left(d_{u}^{2}+d_{1}^{2} m_{u} d_{u}\right) x_{u}^{2}}
$$

Hence

$$
\mu_{1} \leq \max \left\{\sqrt{2\left(d_{u}^{2}+d_{1}^{2} m_{u} d_{u}\right)} \quad: \quad u \in V(G)\right\}
$$

## §3. Lower Bonud for Spectral Radius of Graphs

In this section we establish a lower bound for the spectral radius $\mu_{1}$ of $G$.

Lemma 3.1 ([7][8]) Let $M$ be real symmetric matrix with eigenvalues $\lambda_{1} \geq \lambda_{2} \geq \ldots \geq \lambda_{n}$. Given a partition $\{1,2, \ldots, n\}=\Delta_{1} \cup \Delta_{2} \cup \ldots \cup \Delta_{m}$ with $\left|\Delta_{i}\right|=n_{i}>0$, consider the corresponding blocking $M=\left(M_{i j}\right)$, so that $M_{i j}$ is an $n_{i} \times n_{j}$ block. Let $e_{i j}$ be the sum of the entries in $M_{i j}$ and put $B=\left(\frac{e_{i j}}{n_{i}}\right)$ i.e., $\left(\frac{e_{i j}}{n_{i}}\right.$ is an average row sum in $\left.M_{i j}\right)$. let $\gamma_{1} \geq \gamma_{2} \geq \ldots \geq \gamma_{m}$ be the eigenvalues of $B$. Then the inequalities

$$
\lambda_{i} \geq \gamma_{i} \geq \lambda_{n-m+i} \quad(i=1,2, \ldots, m)
$$

hold. Moreover, if for some integer $k, 1 \leq k \leq m, \lambda_{i}=\gamma_{i}$ for $i=1,2, \ldots, k$ and $\lambda_{n-m+i}=\gamma_{i}$ for $i=k+1, k+2, \ldots, m$, then all the blocks $M_{i j}$ have constant row and column sums.

Let $G$ be a connected graph with $n$ vertices and $m$ edges. Let $V_{1}=\left\{v_{1}, v_{2}, \ldots v_{n_{1}}\right\}$ and $V_{2}=\left\{v_{n_{1}+1}, v_{n_{1}+2} \ldots v_{n}\right\}$ be two partitions of vertices of graph $G$. Let

$$
\begin{array}{cc}
r_{1}=\frac{1}{n_{1}} \sum_{\substack{i, j=1 \\
i \neq j}}^{n_{1}} \min \left(d\left(v_{i}\right), d\left(v_{j}\right)\right), \quad r_{2}=\frac{1}{n-n_{1}} \sum_{\substack{i, j=1 \\
i \neq j}}^{n-n_{1}} \min \left(d\left(v_{n_{1}+i}\right), d\left(v_{n_{1}+j}\right)\right), \\
k_{1}=\frac{-1}{n_{1}} \sum_{\substack{i, j=1 \\
i \neq j}}^{n-n_{1}} \min \left(d\left(v_{i}\right), d\left(v_{n_{1}+j}\right)\right), & k_{2}=\frac{-1}{n-n_{1}} \sum_{\substack{i=1 \\
j=1,2, \ldots, n \\
i \neq j}}^{n-n_{1}} \min \left(d\left(v_{n_{1}+i}\right), d\left(v_{j}\right)\right), \\
d_{1}=\frac{1}{n_{1}} \sum_{v \in V_{1}} d(v), & d_{2}=\frac{1}{n-n_{1}} \sum_{v \in V_{2}} d(v),
\end{array}
$$

where $d(v)$ is the degree of the vertex $v$ of $G$. Now we prove the following theorem.
Theorem 3.2 Let $G$ be a connected graph with $n$ vertices and $m$ edges, then

$$
\mu_{1} \geq \frac{1}{2}\left\{d_{2}+d_{1}-r_{2}-r_{1}+\sqrt{\left(d_{2}-d_{1}-r_{2}+r_{1}\right)^{2}-4 k_{1} k_{2}}\right\} .
$$

Proof Rewrite $L(G)$ as

$$
L(G)=\left(\begin{array}{ll}
L_{11} & L_{12} \\
L_{21} & L_{22}
\end{array}\right)
$$

For $1 \leq i, j \leq 2$, let $e_{i j}$ be the sum of the entries in $L_{i j}$ and put $B=\left(e_{i j} / n_{i}\right)$. Then

$$
B=\left(\begin{array}{cc}
d_{1}-r_{1} & k_{1} \\
k_{2} & d_{2}-r_{2}
\end{array}\right)
$$

and so

$$
|\lambda I-B|=\left|\begin{array}{cc}
\lambda-\left(d_{1}-r_{1}\right) & -k_{1} \\
-k_{2} & \lambda-\left(d_{2}-r_{2}\right)
\end{array}\right| .
$$

Therefore we have

$$
\lambda=\frac{1}{2}\left\{d_{2}+d_{1}-r_{2}-r_{1} \pm \sqrt{\left(d_{2}-d_{1}-r_{2}+r_{1}\right)^{2}-4 k_{1} k_{2}}\right\} .
$$

Thus by Lemma 3.1 we get

$$
\mu_{1} \geq \frac{1}{2}\left\{d_{2}+d_{1}-r_{2}-r_{1}+\sqrt{\left(d_{2}-d_{1}-r_{2}+r_{1}\right)^{2}-4 k_{1} k_{2}}\right\} .
$$

Acknowledgment The first author is thankful to the Department of Science and Technology, Government of India, New Delhi for the financial support under the grant DST/SR/S4/MS: 490107.

## References

[1] C.Adiga and C.S. Shivakumar Swamy, Bounds on the largest of minimum degree eigenvalues of graphs, Int. Math. Forum, $5(37)(2010), 1823-1831$.
[2] N. Alon, Eigenvalues and expanders, Combinatorica, 6(1986), 83-96.
[3] F.R.K.Chung, Eigenvalues of graphs, Proceeding of the International Congress of Mathematicians, Zürich, Switzerland, 1995, 1333-1342.
[4] D. de Caen, An upper bound on the sum of squares of degrees in a graph, Discrete Math., 185(1998), 245-248.
[5] I.Gutman, The energy of a graph, Ber.math.stat.Sekt. Forschungsz.Graz, 103(1978), 1-22.
[6] I.Gutman and Polanski, Mathematical Concepts in Organic Chemistry, Springer-verlag, Berlin, 1986.
[7] W.H.Haemers, Partitioning and Eigenvalues, Eindhoven University of Technology, Memorandum 1976-11; revised version: A generalization of the Higmansims technique, Proc. Kon. Ned. Akad. Wet. A 81(4)(1978), 445-447.
[8] W.H.Haemers, Interlacing eigenvalues and graphs, Linear Algebra and its Applications, 226-228(1995), 593-616.
[9] R.Merris, Laplacian matrices of graphs: a survey, Linear Algebra Appl., 197-198(1994), 143-176.
[10] B.Mohar, Some applications of Laplace eigenvalues of graphs, G. Hahn, G. Sabidussi(Eds.), Graph Symmetry, Kluwer Academic Publishers, Dordrecht, 1997, 225-275.

I want to bring out the secrets of nature and apply them for the happiness of man. I don't know of any better service to offer for the short time we are in the world.

By Thomas Edison, an American inventor.

## Author Information

Submission: Papers only in electronic form are considered for possible publication. Papers prepared in formats, viz., .tex, .dvi, .pdf, or.ps may be submitted electronically to one member of the Editorial Board for consideration in the International Journal of Mathematical Combinatorics (ISSN 1937-1055). An effort is made to publish a paper duly recommended by a referee within a period of 3 months. Articles received are immediately put the referees/members of the Editorial Board for their opinion who generally pass on the same in six week's time or less. In case of clear recommendation for publication, the paper is accommodated in an issue to appear next. Each submitted paper is not returned, hence we advise the authors to keep a copy of their submitted papers for further processing.

Abstract: Authors are requested to provide an abstract of not more than 250 words, latest Mathematics Subject Classification of the American Mathematical Society, Keywords and phrases. Statements of Lemmas, Propositions and Theorems should be set in italics and references should be arranged in alphabetical order by the surname of the first author in the following style:

## Books

[4]Linfan Mao, Combinatorial Geometry with Applications to Field Theory, InfoQuest Press, 2009.
[12]W.S. Massey, Algebraic topology: an introduction, Springer-Verlag, New York 1977.

## Research papers

[6]Linfan Mao, Combinatorial speculation and combinatorial conjecture for mathematics, International J.Math. Combin., Vol.1, 1-19(2007).
[9]Kavita Srivastava, On singular H-closed extensions, Proc. Amer. Math. Soc. (to appear).
Figures: Figures should be drawn by TEXCAD in text directly, or as EPS file. In addition, all figures and tables should be numbered and the appropriate space reserved in the text, with the insertion point clearly indicated.

Copyright: It is assumed that the submitted manuscript has not been published and will not be simultaneously submitted or published elsewhere. By submitting a manuscript, the authors agree that the copyright for their articles is transferred to the publisher, if and when, the paper is accepted for publication. The publisher cannot take the responsibility of any loss of manuscript. Therefore, authors are requested to maintain a copy at their end.

Proofs: One set of galley proofs of a paper will be sent to the author submitting the paper, unless requested otherwise, without the original manuscript, for corrections after the paper is accepted for publication on the basis of the recommendation of referees. Corrections should be restricted to typesetting errors. Authors are advised to check their proofs very carefully before return.

## Contents

Duality Theorems of Multiobjective Generalized Disjunctive Fuzzy Nonlinear Fractional Programming BY E.E.AMMAR ..... 01
Surface Embeddability of Graphs via Joint Trees BY YANPEI LIU ..... 15
Plick Graphs with Crossing Number 1
BY B.BASAVANAGOUD AND V.R.KULLI ..... 21
Effects of Foldings on Free Product of Fundamental Groups BY M.El-GHOUL, A. E.El-AHMADY, H.RAFAT AND M.ABU-SALEEM . ..... 29
Absolutely Harmonious Labeling of Graphs
BY M.SEENIVASAN AND A.LOURDUSAMY ..... 40
The Toroidal Crossing Number of $K_{4, n}$
BY SHENGXIANG LV, TANG LING AND YUANGQIU HUANG ..... 52
On Pathos Semitotal and Total Block Graph of a Tree BY MUDDEBIHAL M. H. AND SYED BABAJAN ..... 64
Varieties of Groupoids and Quasigroups Generated by Linear-Bivariate
Polynomials Over Ring $Z_{n}$ BY E.ILOJIDE, T.G.JAIYÉOLÁ AND O.O.OWOJORI79
New Characterizations for Bertrand Curves in Minkowski 3-Space
BY BAHADDIN BUKCU, MURAT KEMAL KARACAN AND NURAL YUKSEL ..... 98
Respectable Graphs
BY A.SATYANARAYANA REDDY ..... 104
Common Fixed Points for Pairs of Weakly Compatible Mappings BY RAKESH TIWARI AND S.K.SHRIVASTAVA ..... 111
Some Results on Generalized Multi Poly-Bernoulli and Euler Polynomials BY HASSAN JOLANY, HOSSEIN MOHEBBI AND R.EIZADI ALIKELAYE ..... 117
Bounds on the Largest of Minimum Degree Laplician Eigenvalues of a Graph BY CHANDRASHEKAR ADIGA AND C. S. SHIVAKUMAR SWAMY ..... 130

An International Journal on Mathematical Combinatorics


[^0]:    ${ }^{1}$ Received December 09, 2010. Accepted May 8, 2011.

[^1]:    ${ }^{1}$ Received January 28, 2011. Accepted May 12, 2011.

[^2]:    ${ }^{1}$ Received January 6, 2011. Accepted May 18, 2011.

[^3]:    ${ }^{1}$ Received January 1, 2011. Accepted May 20, 2011.

[^4]:    ${ }^{1}$ Received November 08, 2010. Accepted May 22, 2011.

[^5]:    ${ }^{1}$ Supported by National Natural Science Foundation of China (No.10771062) and New Century Excellent Talents in University (No. NCET-07-0276).
    ${ }^{2}$ Received January 8, 2011. Accepted May 25, 2011.

[^6]:    ${ }^{1}$ Received November 23, 2010. Accepted May 26, 2011.

[^7]:    ${ }^{1}$ Received February 17, 2011. Accepted May 28, 2011.

[^8]:    ${ }^{1}$ Received October 29, 2010. Accepted May 30, 2011.

[^9]:    ${ }^{1}$ Received February 23, 2011. Accepted June 2, 2011.

[^10]:    ${ }^{1}$ Received October 20, 2010. Accepted June 6, 2011.

[^11]:    ${ }^{1}$ Received December 29, 2010. Accepted June 8, 2011.

[^12]:    ${ }^{1}$ Received November 3, 2010. Accepted June 10, 2011.

