# Hamiltonian Cycles in Symmetric Graphs 

## Dov Zazkis

B.Sc., Simon Fraser University, 2006

Thesis Submitted In Partial Fulfillment Of
The Requirements For The Degree Of
Master of Science
in
Mathematical, Computer, and Physical Sciences
(Mathematics)

University of Northern British Columbia
November 2008
(C) Dov Zazkis, 2008

Library and
Archives Canada
Published Heritage Branch

395 Wellington Street Ottawa ON K1A 0N4 Canada

Bibliothèque et
Archives Canada
Direction du
Patrimoine de l'édition
395, rue Wellington Ottawa ON K1Ă 0N4
Canada

Your file Votre référence
ISBN: 978-0-494-48775-4
Our file Notre référence
ISBN: 978-0-494-48775-4

## NOTICE:

The author has granted a nonexclusive license allowing Library and Archives Canada to reproduce, publish, archive, preserve, conserve, communicate to the public by telecommunication or on the Internet, loan, distribute and sell theses worldwide, for commercial or noncommercial purposes, in microform, paper, electronic and/or any other formats.

The author retains copyright ownership and moral rights in this thesis. Neither the thesis nor substantial extracts from it may be printed or otherwise reproduced without the author's permission.

## AVIS:

L'auteur a accordé une licence non exclusive permettant à la Bibliothèque et Archives Canada de reproduire, publier, archiver, sauvegarder, conserver, transmettre au public par télécommunication ou par l'Internet, prêter, distribuer et vendre des thèses partout dans le monde, à des fins commerciales ou autres, sur support microforme, papier, électronique et/ou autres formats.

L'auteur conserve la propriété du droit d'auteur et des droits moraux qui protège cette thèse. Ni la thèse ni des extraits substantiels de celle-ci ne doivent être imprimés ou autrement reproduits sans son autorisation.

In compliance with the Canadian Privacy Act some supporting forms may have been removed from this thesis.

While these forms may be included in the document page count, their removal does not represent any loss of content from the thesis.

Conformément à la loi canadienne sur la protection de la vie privée, quelques formulaires secondaires ont été enlevés de cette thèse.

Bien que ces formulaires aient inclus dans la pagination, il n'y aura aucun contenu manquant.


#### Abstract

Let $k$ be a positive integer. We define $M_{k}$ to be the graph with a vertex set consisting of all binary strings of length $2 k+1$ which have either $k$ or $k+1$ ones and edge set consisting of all pairs of these binary strings which differ in exactly one bit. Showing that the graph $M_{k}$ is Hamiltonian for all $k$ is known as the Middle Levels problem. This problem was first posed in the early 1980's and to this day remains unsolved. In this thesis we explore the symmetries of $M_{k}$ and graphs related to it. We then use these symmetries to propose a method for finding Hamiltonian cycles in $M_{k}$ when $2 k+1$ and $k$ are prime. We believe that our method is more efficient than methods proposed by previous authors..


## Acknowledgements

I'd like to thank my supervisor, Dr. Iliya Bluskov, for his guidance, support and for introducing me to the Middle Levels problem which I have enjoyed working on so much. I'd like to thank Dr. Jennifer Hyndman, for helping me negotiate through some of the group theory problems which unexpectedly popped up during the writing of this thesis. I'd like to thank my mother, Dr. Rina Zazkis, for her endless emotional and almost endless financial support. Finally, I'd like to thank Heinrich Butow, Eric Chlebek and Ashton Fedler for helping me with coding of my algorithm.

## Contents

1 Introduction ..... 1
2 Density Based Results ..... 4
2.1 Definitions ..... 4
2.2 Some Well Known Results ..... 7
3 Hamiltonian Cycles in Symmetric Graphs ..... 10
3.2 The Middle Levels Problem ..... 12
3.3 Definitions ..... 14
3.4 Properties of $M_{k}$ ..... 16
3.5 Properties of $R_{k}$ ..... 33
3.6 The Automorphisms of $R_{k}$ ..... 44
3.7 Catalan Numbers ..... 46
4 Reducing the Middle Levels Problem ..... 49
4.1 A Reduction of the Middle Levels Problem ..... 49
4.2 Another Possible Reduction ..... 54
4.3 Implementation ..... 59
4.4 Some Possible Improvements the Algorithm ..... 60
5 Applications to Gray Codes ..... 61
6 Conclusion ..... 66
Bibliography ..... 67
Appendix ..... 69
A Java code ..... 70
B Code output ..... 79
C Diagrams of $M_{k}$ 's and $R_{k}$ 's ..... 80
C. $1 \quad R_{k}$ ..... 80
C. $2 M_{k}$ ..... 82

## List of Figures

3.1 All known non-Hamiltonian vertex transitive graphs ..... 11
3.2 Midlevel Hotel ..... 13
3.3 Déjter's reduction 1 ..... 34
3.4 Déjter's reduction 2 ..... 35
3.5 Déjter's reduction 3 ..... 35
$3.6 R_{3}$ colored path ..... 36
3.7 $M_{3}$ with colored Hamiltonian cycle ..... 37
$3.8 \quad M_{3}$ explicit colored Hamiltonian cycle ..... 38
3.9 CFA's of $R_{3}$ ..... 46
$4.1 \quad R_{5}$ ..... 51
$4.2 \quad Z_{5}$ subgraphs of $R_{5}$ ..... 52
$4.3 \quad Z_{5}$ cycles of $R_{5}$ ..... 53
4.4 Hamiltonian path in $R_{5}$ ..... 54
4.5 Pósa path heuristic ..... 56
4.6 Modified Pósa path heuristic ..... 58
5.1 Gray code spinning disk . . . . . . . . . . . . . . . . . . . . . . . . . 62

## List of Tables

3.1 Anti-Palindromes are not $l$-vertices ..... 23
5.1 Binary Counting ..... 62
5.2 Gray Code Counting ..... 63

## Chapter 1

## Introduction

A graph $G$ is defined as a pair of sets $(V(G), E(G))$, where $V(G)$ is a set of objects called vertices and $E(G)$ is a set of unordered pairs of vertices from $V(G)$ called edges. A Hamiltonian cycle in $G$ is a list of vertices $v_{1}, v_{2}, v_{3}, \ldots v_{n}$ such that every vertex in $V(G)$ appears in the list exactly once and $v_{i} v_{i+1}, v_{1} v_{n} \in E(G)$ for $1 \leq i<n$. Studying the existence of Hamiltonian cycles in graphs has been an important part of graph theory since its inception. Although graphs were not explicitly mentioned until 1878, some examples of problems equivalent to finding Hamiltonian cycles in certain graphs date back hundreds of years. The knights tour problem is such an example. It involves finding a sequence of moves on a chessboard which allows a knight to visit each of the 64 squares on that chessboard exactly once and return to the starting point.


This problem can be represented as a graph: Let each vertex represent a square on the chessboard and for each pair of vertices, $v_{1}, v_{2}$, let $v_{1} v_{2}$ be an edge in the graph if and only if the squares which $v_{1}$ and $v_{2}$ represent can be reached from one another by a knight in one move. A Hamiltonian cycle in this graph is a solution to the knights tour problem. Systematic solutions to this problem were published by Leonhard Euler in 1759 [10].

The studies of Hamiltonian cycles in graphs can be roughly split into two main categories of problems: 1) Finding sets of conditions which force the existence of sufficiently many edges which in turn forces the existence of a Hamiltonian cycle, 2) Using symmetries as well as structural characteristics to show that certain classes of graphs have Hamiltonian cycles. This is typically done recursively, using Hamiltonian cycles in smaller graphs to build Hamiltonian cycles in larger graphs. For convenience, we will refer to the first type as "density based" and the second type as "non-density based" results. Our main focus will be the latter.

Density based theorems are limited in that they can not be used to show that certain
classes of graphs which have a relatively low number of edges have a Hamiltonian cycle. For example, none of the density based theorems can be used to show that $C_{k}$ with $k>6$ ( the cycle on $k$ vertices) is Hamiltonian, even though it is trivial to show that $C_{k}$ is Hamiltonian. There are many classes of graphs which are believed to be Hamiltonian for which density based theorems can not be used to show the existence of a Hamiltonian cycle. These graphs are typically quite symmetric. To show that graphs which are believed to be Hamiltonian are in fact Hamiltonian, it becomes necessary to use the properties of these graphs to either explicitly construct a Hamiltonian cycle or show its existence using other non-density based methods. These non-density based methods are also limited in that all current known results are difficult to apply to a wide variety of graphs. There are many open problems which involve proving that certain classes of highly symmetric graphs are Hamiltonian.

## Chapter 2

## Density Based Results

Even though it is not our focus, it is reasonable to start with some density based results to give the reader a sense of what kind of results can be applied generally to all graphs to test whether they are Hamiltonian, regardless of symmetries. We will also include some well known group theory results in this section which will be needed in later proofs.

### 2.1 Definitions

Definition 2.1.1 A graph, $G$, is a pair $(V(G), E(G))$, where $V(G)$ is a set of points called vertices and $E(G)$ is set of unordered pairs of vertices from $V(G)$.

Definition 2.1.2 Two vertices, $v_{1}$ and $v_{2}$ in $V(G)$ are said to be adjacent if $v_{1} v_{2} \in$ $E(G)$.

Definition 2.1.3 A subgraph, $G^{\prime}=\left(V^{\prime}(G), E^{\prime}(G)\right)$, of a graph $G=(V(G), E(G))$ is a graph such that $V^{\prime}(G) \subseteq V(G)$ and $E^{\prime}(G) \subseteq E(G)$ with $x y \in E^{\prime}(G)$ only if $x, y \in V^{\prime}(G)$. The subgraph is said to be a spanning subgraph if $V^{\prime}(G)=V(G)$.

Definition 2.1.4 The subgraph of $G=(V(G), E(G))$ induced by the set $W \subseteq$ $V(G)$ is the subgraph $G^{\prime}=\left(V^{\prime}(G), E^{\prime}(G)\right)$ with $V^{\prime}(G)=W$ and $E^{\prime}(G)=\left\{v_{i} v_{j} \mid v_{i}, v_{j} \in\right.$ $V^{\prime}(G)$ and $\left.v_{i} v_{j} \in E(G)\right\}$.

Definition 2.1.5 A graph $G=(V(G), E(G))$ is said to be bipartitie if $V(G)=$ $X \cup Y$ with $|X \cap Y|=0$, and if $v_{1} v_{2} \in E(G)$, then exactly one of $v_{1}, v_{2}$ is in $X$.

Definition 2.1.6 A graph $G$ is said to be Hamiltonian if it contains a Hamiltonian cycle.

Definition 2.1.7 A Hamiltonian path in a graph $G=(V(G), E(G))$ is a list of vertices $v_{1}, v_{2}, v_{3}, \ldots v_{n}$ such that every vertex in $V(G)$ appears exactly once in the list and $v_{i} v_{i+1} \in E(G)$ for $1 \leq i<n$.

Definition 2.1.8 A path in a graph $G=(V(G), E(G))$ is a list $v_{1}, v_{2}, v_{3}, \ldots, v_{k}$ such that $v_{i} v_{i+1} \in E(G)$ for $1 \leq i<k$ and for all $1 \leq i, j<k v_{i} \neq v_{j}$. The vertices $v_{1}$ and $v_{k}$ are referred to as the endpoints of the path.

Definition 2.1.9 $A$ graph $G$ is said to be $k$-connected if for every pair of vertices $v_{1}, v_{2}$ the graph contains at least $k$ disjoint paths with $v_{1}, v_{2}$ as endpoints. In general, we refer to 1-connected graphs as connected graphs.

Definition 2.1.10 We say a set $X \subseteq V(G)$ is independent if there are no edges between vertices in $X$. The cardinality of the largest independent set in $G$ is called
the independence number and is denoted $\alpha(G)$.
Definition 2.1.11 The degree of a vertex $v$, denoted $\delta(v)$ is the number of edges incident with it. The minimum and maximum degree of a graph $G$ are defined as $\delta_{\min }(G)=\min _{v \in V(G)} \delta(v)$ and $\delta_{\max }(G)=\max _{v \in V(G)} \delta(v)$, respectively.

Definition 2.1.12 A graph $G$ is said to be m-regular if $\delta(v)=m$ for all $v \in V(G)$.
Definition 2.1.13 Two graphs, $G=(V(G), E(G))$ and $H=(V(H), E(H))$ are said to be isomorphic if there exists a bijective mapping $\gamma: V(G) \mapsto V(H)$ such that if $v_{1}, v_{2} \in V(G)$ and $\gamma\left(v_{1}\right)=v_{1}^{\prime}, \gamma\left(v_{2}\right)=v_{2}^{\prime}$ then $v_{1}$ and $v_{2}$ are adjacent if and only if $v_{1}^{\prime}$ and $v_{2}^{\prime}$ are adjacent.

Definition 2.1.14 Two graphs, $G=(V(G), E(G))$ and $H=(V(H), E(H))$, are said to be homomorphic if there exists a mapping $\gamma: V(G) \mapsto V(H)$ such that if $v_{1}, v_{2} \in V(G)$ and $\gamma\left(v_{1}\right)=v_{1}^{\prime} \neq v_{2}^{\prime}=\gamma\left(v_{2}\right)$ then $v_{1}$ and $v_{2}$ are adjacent if and only if $v_{1}^{\prime}$ and $v_{2}^{\prime}$ are adjacent.

Definition 2.1.15 A automorphism $\phi$ of a graph $G$ is a isomorphic mapping from $G$ to itself.

Definition 2.1.16 $A$ graph $G$ is symmetric if it has a non-trivial automorphism.
Definition 2.1.17 The Euler $\phi$ function of a positive integer n, denoted $\phi(n)$, is the number of positive integers $i$ less than $n$ such that $i$ and $n$ are relatively prime.

Definition 2.1.18 A group is a set $\mathbb{G}$ together with a binary operation $\otimes$ that satisfies the 4 properties:
(a) $\mathbb{G}$ is closed under the operation $\otimes$, meaning that if $x, y$ are in $\mathbb{G}$ then so is $x \otimes y$.
(b) The operation $\otimes$ is associative, that is, $(x \otimes y) \otimes z=x \otimes(y \otimes z)$.
(c) There is an element in $\mathbb{G}$ called e such that $e \otimes x=x=x \otimes e$ for all $x \in \mathbb{G}$.
(d) Each element $x$ of $\mathbb{G}$ has a so-called inverse, $x^{-1}$, such that $x \otimes x^{-1}=e=x^{-1} \otimes x$.

Note 2.1.19 If in addition a group has the property that $x \otimes y=y \otimes x$ for all $x, y \in \mathbb{G}$, then the group is called an Abelian group. We will only deal with finite Abelian groups in this thesis.

Definition 2.1.20 A subgroup $\mathbb{H}$ of a group $\mathbb{G}$ is a non-empty subset of $\mathbb{G}$ which is a group under the same binary operation as $\mathbb{G}$.

Definition 2.1.21 A coset of a subgroup $\mathbb{H}$ of a group $\mathbb{G}$, denoted $g \otimes \mathbb{H}$, is the set $\{g \otimes h \mid h \in \mathbb{H}\}$.

Definition 2.1.22 The subgroup generated by $x \in \mathbb{G}$ is the set $e, x, x \otimes x, x \otimes$ $x \otimes x, \ldots, x^{-1}$. This set is denoted $\langle x\rangle$.

Note 2.1.23 $\langle x\rangle$ is in fact a subgroup of $\mathbb{G}$.
Definition 2.1.24 The order of a group $\mathbb{G}$ is the cardinality of the set of its elements.

### 2.2 Some Well Known Results

We will list some theorems that are useful in deciding whether a dense graph is Hamiltonian.

THEOREM 2.2.1 [18] If $\delta(v)+\delta(u) \geq|V(G)|$ for every pair of distinct nonadjacent vertices $u, v \in V(G)$, then $G$ is Hamiltonian.

One consequence of Theorem 2.2.1 the following corollary:
Corollary 2.2.2 [8] If $G$ is a graph such that $\delta_{\min }(G) \geq|V(G)| / 2$, then $G$ is Hamiltonian.

THEOREM 2.2.3 [16] If $G=(X \cup Y, E)$ is a bipartite graph with $|X|=n=|Y|$ ( $n \geq 2$ ) and $\delta(u)+\delta(v) \geq n+1$ for each non-adjacent pair $u \in X$ and $v \in Y$, then $G$ is Hamiltonian.

THEOREM 2.2.4 [14] If $G$ is a d-regular 2-connected graph of order $n$ with $d \geq$ $|V(G)| / 3$, then $G$ is Hamiltonian.

The following are well-known group theory results which we will utilize later.
THEOREM 2.2.5 Let $\mathbb{H}$ be a subgroup of the group $\mathbb{G}$. Then any two cosets $g_{1} \otimes \mathbb{H}$, $g_{2} \otimes \mathbb{H}$ are either disjoint or identical and the order of any coset is equal to the order of the subgroup used to generate it.

THEOREM 2.2.6 (Lagrange's Theorem) If $\mathbb{H}$ is a subgroup of $\mathbb{G}$ then $|\mathbb{H}|$ divides $|\mathbb{G}|$.

THEOREM 2.2.7 (First Sylow Theorem) Let $p$ be prime and let $\mathbb{H}$ be a group. If $p^{m}$ is the highest power of $p$ which is a factor of $|\mathbb{H}|$, then $\mathbb{H}$ has a subgroup of order $p^{m}$.

THEOREM 2.2.8 (Third Sylow Theorem) Let $p, m$ and $\mathbb{H}$ be the same as in the
pervious theorem. Then the number of subgroups of $\mathbb{H}$, which have order $p^{m}$ is congruent to 1 modulo $p$.

## Chapter 3

## Hamiltonian Cycles in Symmetric

## Graphs

Graphs are often represented as "connect the dots" diagrams were each dot represents a vertex and two dots are connected by a line if the two vertices they represent are adjacent. Because we define graphs in terms of sets and often work with graphs in terms of dots and lines they tend to be thought of as both set theoretic and geometric objects. This sometimes leads to two words, one set theoretic and one geometric in origin being used to describe the same thing. For example; symmetries and automorphisms are two words we use to describe the same thing. A graph $G$ is said to be vertex-transitive if for every distinct pair of vertices $v_{1}, v_{2} \in V(G)$ there exists an automorphism of $G$ which maps $v_{1}$ to $v_{2}$. We can intuitively think of a symmetry of a graph as a rearrangement of that graph's vertices which preserves
edges just like in a geometric setting. A vertex-transitive graph can then be thought of as "as symmetric as possible", because each pair of vertices can be thought of as "the same" under one of the symmetries (automorphisms) of that graph. It is known that not all vertex-transitive graphs are Hamiltonian. There are 6 known non-Hamiltonian connected vertex-transitive graphs (pictured below).


Figure 3.1:

Arguably, the most important conjecture relating to Hamiltonian Cycles in symmetric graphs is the following, due to Lovász [17] .

Conjecture 3.1.9 a) Every connected vertex-transitive graph has a Hamiltonian path.
b) Aside from the six known counterexamples, all connected vertex-transitive graphs
are Hamiltonian.

Part b) of conjecture 3.1.9 encompasses the work done in this thesis since the graphs we will examine are vertex-transitive.

### 3.2 The Middle Levels Problem

One problem which Lovász's conjecture encompasses is known as the Middle Levels problem. It was first proposed in a paper of Havel's in 1982 [12] although various authors have falsely attributed the problem to authors other than Havel (most notably to: Edrös, Kelly, Déjter, Trotter). The following is Edrös's proposal of the problem:

Suppose that there is a hotel with $k+1$ people in the lobby and $k$ people outside. There is a single door which leads from the outside to the inside and only one person can enter or exit at a time. There are $2\binom{2 k+1}{k}$ ways of placing each of the $2 k+1$ people either in or out of the hotel such that either $k$ or $k+1$ end up inside. Is it always possible to go through all of these arrangements exactly once and end up back at the start by letting one person at a time in or out of the hotel?


Figure 3.2:

The above interpretation of the Middle Levels problem is useful for explaining the problem, however, this formulation is difficult to use in analyzing the problem. So let us change how the problem is posed. Assign the $2 k+1$ people some arbitrary order. Assign the $i^{\text {th }}$ person the $i^{\text {th }}$ bit in a $2 k+1$ bit string and make that bit 1 if $i^{\text {th }}$ person is inside the hotel and 0 otherwise. It is now easy to see that our problem is equivalent to the following:

For a fixed $k$, let $X$ be the set of all $2 k+1$ bit binary strings containing exactly $k$ 1's and let $Y$ be the set of all $2 k+1$ bit binary strings containing exactly $k+1$ 1's. Now let $M_{k}$ be the the graph with vertex set $X \cup Y$ and $x \in X$ adjacent to $y \in Y$ if and only if $x$ differs from $y$ in exactly one bit. Is $M_{k}$ Hamiltonian for all $k$ ?

This problem has become known as the "Middle Levels Problem" because $M_{k}$ can be thought of as the subgraph induced by the two middle levels of the boolean lattice.

It is conjectured that $M_{k}$ is Hamiltonian for all $k>0$. Savage, Shields and Shields [25] showed that $M_{k}$ is Hamiltonian for $1 \leq k \leq 17$. Although many people have worked on this problem, there have been no published thorough explorations of the properties of $M_{k}$ and related graphs. We will explore some of these properties and use them to propose improvements to the algorithm used in [25]. Our focus will be the case where $2 k+1$ and $k$ are both prime.

### 3.3 Definitions

For convenience, we will set $n=2 k+1$ for the remainder of this thesis. Also we will take all subscripts and superscripts to be modulo $n$.

Definition 3.3.1 Given $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle \in V\left(M_{k}\right)$ we call the circular shift of that vertex, denoted $\operatorname{sh}(A)$, the vertex $\left\langle a_{n}, a_{1}, a_{2}, \ldots, a_{n-1}\right\rangle$. The $j^{\text {th }}$ circular shift of $A$ is denoted $\operatorname{sh}^{j}(A)=\underbrace{\operatorname{sh}(\operatorname{sh}(\ldots \ldots . \operatorname{sh}}_{j \text { times }}(A))$. We also define $\operatorname{sh}^{0}(A)=A$.

Note 3.3.2 If we are given an $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle \in M_{k}$, then

$$
s h^{j}(A)=\left\langle a_{1-j}, a_{2-j}, \ldots, a_{n-j}\right\rangle
$$

Definition 3.3.3 Given $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle \in V\left(M_{k}\right)$, we call the flip of that vertex, denoted $f l(A)$, the vertex $\left\langle a_{n}, a_{n-1}, a_{n-2}, \ldots, a_{1}\right\rangle$.

Definition 3.3.4 The class of $A \in V\left(M_{k}\right)$ is the set $\left\{B \in V\left(M_{k}\right) \mid B=s h^{j}(A), 0 \leq\right.$
$j \leq n-1\}$. If $A, B \in V\left(M_{k}\right)$ with $\operatorname{sh}^{m}(A)=B$ for some $m$, then we say that $B$ is in the class of $A$, denoted $B \in \operatorname{cl}(A)$, or, equivalently, that $A$ and $B$ are in the same class.

Definition 3.3.5 Given $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle \in V\left(M_{k}\right)$, we call the complement of $A$, denoted $A^{c}$, the vertex $\left\langle 1-a_{1}, 1-a_{2}, \ldots, 1-a_{n}\right\rangle$.

Note 3.3.6 Each class has exactly $n$ vertices in it, they are disjoint, and every class has a unique complement class.

Definition 3.3.7 We call a vertex $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle \in V\left(M_{k}\right)$ an l-vertex if $A$ is adjacent to a vertex in the class of $A^{c}$.

Definition 3.3.8 We call cl $(A)$ an l-class if $A$ is an l-vertex. See Remark 3.4.3 for more details.

Example 3.3.9 1010100 is an l-vertex because it is adjacent to $\operatorname{sh}\left((1010100)^{c}\right)=$ 1010101.

Definition 3.3.10 $A$ bit string, $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$, is a palindrome if $f l(A)=A$.

Example 3.3.11 11001010011 is a palindrome.
Definition 3.3.12 $A$ bit string, $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$, is an anti-palindrome if $f l\left(A^{c}\right)$ differs from $A$ only in the $(k+1)^{\text {th }}$ position.

Example 3.3.13 110100100 is an anti-palindrome.

Note 3.3.14 The concepts 'class' and 'l-vertices' were used in previous work on the middle levels problem, specifically [4], [5], [22], [23], [24] and [25], but were given
no formal definitions or names. However, these concepts were utilized inside proofs. Since we will be referring to these concepts throughout this thesis it was necessary to formalize them.

### 3.4 Properties of $M_{k}$



The above four diagrams represent the first four $M_{k}$ 's. The graph on the left, $M_{1}$, is isomorphic to $C_{6}$, so the graph is itself a cycle and hence clearly Hamiltonian. However, the number of vertices in $M_{k}$ grows exponentially and finding a Hamiltonian cycle becomes progressively harder.

Observation 3.4.1 If a vertex $A$ in $M_{k}$ is adjacent to $B$, then $\operatorname{sh}(A)$ is adjacent to $\operatorname{sh}(B)$ and $A^{c}$ is adjacent to $B^{c}$. Furthermore, if $B_{i} \in c l\left(B_{j}\right)$, then $B_{i}^{c} \in \operatorname{cl}\left(B_{j}^{c}\right)$ and $\operatorname{sh}\left(B_{i}\right) \in \operatorname{cl}\left(\operatorname{sh}\left(B_{j}\right)\right)$.

Lemma 3.4.2 If $A$ is adjacent to $B_{1}$ and $B_{2}$ from the same class, then there exist $C \neq A$ such that $C$ is in the same class as $A$ and $B_{1}$ is adjacent to both $A$ and $C$.

Proof: Since $B_{1}$ and $B_{2}$ are in the same class, $\operatorname{sh}^{j}\left(B_{2}\right)=B_{1}$ for some $1 \leq j<n$. From Observation 3.4.1, we know that $s h^{j}(A)$ is adjacent to $s h^{j}\left(B_{1}\right)$ and $s h^{j}\left(B_{2}\right)$.

But $s h^{j}\left(B_{2}\right)=B_{1}$, which means that $B_{1}$ is adjacent to $s h^{j}(A)$ and A. Hence we can take $C=s h^{j}(A)$, which completes the proof. (Note that since $1 \leq j<n$, $\left.s h^{j}(A) \neq A\right)$.

Observation 3.4.3 From Observation 3.4.1 we can conclude that if one vertex in a class is an l-vertex, then all the vertices in that class must be l-vertices.

THEOREM 3.4.4 $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$ is an l-vertex if and only if there exist $j$ and $s$ such that $a_{i}=1-a_{i-j}$ for all $i \in\{1,2, . . n\} \backslash\{s\}$ and $a_{s}=a_{s-j}$.

Proof: Assume $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$ is an $l$-vertex. This means that there exist $j$, $1 \leq j<n$, such that the vertex $\operatorname{sh}^{j}\left(A^{c}\right)=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, \ldots, a_{n}^{\prime}\right\rangle$ differs from $A$ in exactly one position. Let that position be $s$. Then $a_{s} \neq a_{s}^{\prime}=1-a_{s-j}$, which implies $a_{s}=a_{s-j}$. We also have $a_{i}=a_{i}^{\prime}=1-a_{i-j}$ for $i \neq s$. So A has the desired property.

Now, assume instead that $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$ has the properties $a_{i}=1-a_{i-j}$ for $i \neq s$ and $a_{s}=a_{s-j}$ for some $j$ and $s$. Without loss of generality, we assume that $s=1$. This means that $s h^{j}\left(A^{c}\right)=\left\langle 1-a_{1-j}, 1-a_{2-j}, \ldots, 1-a_{n-j}\right\rangle=\left\langle 1-a_{1}, a_{2}, a_{3}, a_{4}, \ldots a_{n}\right\rangle$. Hence $A$ differs from $s h^{j}\left(A^{c}\right)$ in only the first position, and therefore, A is an $l$-vertex.

Lemma 3.4.5 Let $A=\left\langle a_{1}, a_{2}, \ldots a_{n}\right\rangle$ be an l-vertex with $j$ and $s$ as in Theorem 3.4.4, that is, $a_{i}=1-a_{i-j}$ for all $i \in\{1,2, . ., n\} \backslash\{s\}$ and $a_{s}=a_{s-j}$. Then $\operatorname{gcd}(n, j)=1$.

Proof: Suppose toward a contradiction that this is not the case and $\operatorname{gcd}(n, j)>1 . \mathbb{Z}_{n}$
is a group under addition. Since $g c d(n, j) \neq 1$, the order of the subgroup generated by $j$ is less than the order of the group, which is $n$. Lagrange's Theorem tells us that the order of a subgroup divides the order of the group. Since $n$ is odd every divisor of $n$ must also be odd so $|\langle j\rangle|$ must also be odd. Let $x+\langle j\rangle$ be a coset of $\langle j\rangle$ which does not contain $s$. The distinct elements of $x+\langle j\rangle$ form the set $\{x, x+j, x+2 j, \ldots x+(|\langle j\rangle|-1) j\}$ which is the same as $\{x, x-j, x-2 j, \ldots x-(|\langle j\rangle|-1) j\}$ modulo $n$. Note that $x-(|\langle j\rangle|) j \equiv_{n} x$. Since $s \notin x+\langle j\rangle$, Theorem 3.4.4 gives $a_{x}=1-a_{x-j}=a_{x-2 j}=\ldots=a_{x-(\mid\langle j|-1) j}=1-a_{x-(\mid\langle j\rangle) j}=1-a_{x}$, which is a contradiction since $a_{x} \neq 1-a_{x}$.

Observation 3.4.6 It was shown in the proof of Theorem 3.4.4 that if

$$
A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle
$$

is an l-vertex then there is a pair $s, j$ such that $a_{i}=1-a_{i-j}$ for all $i \neq s$. If we let $s^{\prime}=s-j$, then we have $a_{i}=1-a_{i-j}$ for all $i \neq s^{\prime}+j$. This gives us the $n$ equations:

$$
\begin{aligned}
& a_{1}=1-a_{1-j} \\
& a_{2}=1-a_{2-j} \\
& \vdots \\
& a_{s^{\prime}}=1-a_{s^{\prime}-j} \\
& \vdots \\
& a_{s^{\prime}+j-1}=1-a_{s^{\prime}-1} \\
& a_{s^{\prime}+j}=a_{s^{\prime}} \\
& a_{s^{\prime}+j+1}=1-a_{s^{\prime}+1} \\
& \vdots \\
& a_{s^{\prime}+2 j}=1-a_{s^{\prime}+j} \\
& \vdots \\
& a_{n}=1-a_{n-j} .
\end{aligned}
$$

Now, the integers $0,1,2, . ., n-1$ form a complete system of residues modulo $n$. Since Lemma 3.4.5 tells us that $\operatorname{gcd}(n, j)=1$, the integers $s^{\prime}, s^{\prime}-j, s^{\prime}-2 j, . ., s^{\prime}-(n-1) j$ also form a complete system of residues modulo $n$. Therefore, $\left\{s^{\prime}, s^{\prime}-j, s^{\prime}-2 j, . ., s^{\prime}-\right.$ $(n-1) j\} \equiv_{n}\{0,1,2, . ., n-1\}$. Also, observe that $s^{\prime}-(n-1) j \equiv s^{\prime}+j(\bmod n)$, so we can write $a_{s^{\prime}}=1-a_{s^{\prime}-j}=a_{s^{\prime}-2 j}=1-a_{s^{\prime}-3 j}=\ldots=a_{s^{\prime}-(n-1) j}=a_{s^{\prime}+j}$. Since $s^{\prime}, s^{\prime}-j, s^{\prime}-2 j, . ., s^{\prime}-(n-1) j$ is a complete system of residues modulo $n$, all the bits of $A$ are represented in the above string of equalities. In addition, if we know whether $a_{s^{\prime}}$ is equal to 0 or 1 , then we can use the values $s^{\prime}, j$ and $a_{s^{\prime}}$ to construct the bitstring that represents $A$. So $A$ is completely determined by the triple $\left(s^{\prime}, j, a_{s^{\prime}}\right)$. Now, notice that the vertex determined by the triple $\left(s^{\prime}+1, j, a_{s^{\prime}+1}=a_{s^{\prime}}\right)$ is $s h(A)$. This means that all the vertices in the same class as $A$ are determined by $\left(q+s^{\prime}, j, a_{q+s^{\prime}}=a_{s^{\prime}}\right)$
with $1 \leq q \leq n$. Also, the complement class can be determined by triples of the form $\left(q+s^{\prime}, j, a_{q+s^{\prime}}=1-a_{s^{\prime}}\right)$ with $1 \leq q \leq n$.

Lemma 3.4.7 If $A$ is an l-vertex, then some vertex in $c l(A)$ is a palindrome.
Proof: Assume that $A$ is an $l$-vertex determined by the triple $\left(s, j, a_{s}\right)$. Then

$$
a_{s}=1-a_{s-j}=a_{s-2 j}=1-a_{s-3 j}=\ldots \ldots=a_{s+3 j}=1-a_{s+2 j}=a_{s+j},
$$

which gives

$$
\begin{gathered}
a_{s}=a_{s+j} \\
a_{s-j}=a_{s+2 j} \\
a_{s-2 j}=a_{s+3 j} \\
\vdots \\
a_{s-(k-1) j}=a_{s+k j} \\
a_{s-k j}=a_{s+(k+1) j}
\end{gathered}
$$

Note that $s+(k+1) j \equiv_{n} s-k j$. Now we can rearrange the above equations to get:

$$
\begin{aligned}
& a_{s+k j}=a_{s+k j} \\
& a_{s+(k-1) j}=a_{s+(k+1) j} \\
& a_{s+(k-2) j}=a_{s+(k+2) j} \\
& \vdots \\
& a_{s}=a_{s+j}
\end{aligned}
$$

Observe that, $a_{s+(k-i) j}=a_{s+(k+i) j}$ for all $1 \leq i \leq k$. There exists some circular shift of $A^{\prime}$ of $A$ such that the $k^{\text {th }}$ bit of $A^{\prime}$ is $a_{s+k j}$. The $(k+i j)^{t h}$ and $(k-i j)^{\text {th }}$ bits of $A^{\prime}$ are $a_{s+(k+i) j}$ and $a_{s+(k-i) j}$, respectively. Since $k-(i) j \equiv_{n} n-(k+(i) j), A^{\prime}$ is a palindrome, as desired.

Lemma 3.4.8 If $A$ is an l-vertex, then some vertex in $c l(A)$ is an anti-palindrome.

Proof: Let $A$ be an $l$-vertex determined by the triple $\left(s, j, a_{s}\right)$. The vertex $B=$ $\left\langle b_{1}, b_{2}, . ., b_{n}\right\rangle$, determined by the triple $\left(k+1, j, b_{k+1}=a_{s}\right)$, is in $c l(A)$. We have:

$$
b_{k+1}=1-b_{k+1-j}=b_{k+1-2 j}=1-b_{k+1-3 j}=\ldots \ldots=b_{k+1+3 j}=1-b_{k+1+2 j}=b_{k+1+j} .
$$

This gives us:

$$
\begin{gathered}
b_{k+1}=b_{k+1} \\
1-b_{k+1-j}=b_{k+1+j} \\
1-b_{k+1-2 j}=b_{k+1+2 j} \\
1-b_{k+1-3 j}=b_{k+1+3 j} \\
\vdots \\
1-b_{k+1-(k) j}=b_{k+1+k j} .
\end{gathered}
$$

Hence $1-b_{k+1-i j}=b_{k+1+i j}$ for all $1 \leq i \leq k$. Notice that $k+1-i j \equiv_{n} n-(k+1+i j)$. This means that $1-b_{-i}=b_{i}$ when $i \neq k+1$, which implies that B is an antipalindrome, as desired.

Observation 3.4.9 It may seem that Lemmas 3.4.8 and 3.4.7 provide two more useful characterizations of l-vertices, however, since both these proofs are not 'if and only if' proofs it might happen that some non-l-vertices are palindromes or antipalindromes. A simple example will suffice to show that this, in fact, is the case. It is easy to see that $A=\langle 1,1,0,0,1,0,1,1,0,1,1,0,0\rangle$ is an anti-palindrome. The following are all the bit strings in the inverse class of $A$ :

| string in $c l\left(A^{c}\right)$ | bits different from the bits of $A$ |
| :---: | :---: |
| $\langle 0,0,1,1,0,1,0,0,1,0,0,1,1\rangle$ | 13 |
| $\langle 1,0,0,1,1,0,1,0,0,1,0,0,1\rangle$ | 5 |
| $\langle 1,1,0,0,1,1,0,1,0,0,1,0,0\rangle$ | 2 |
| $\langle 0,1,1,0,0,1,1,0,1,0,0,1,0\rangle$ | 9 |
| $\langle 0,0,1,1,0,0,1,1,0,1,0,0,1\rangle$ | 7 |
| $\langle 1,0,0,1,1,0,0,1,1,0,1,0,0\rangle$ | 5 |
| $\langle 0,1,0,0,1,1,0,0,1,1,0,1,0\rangle$ | 7 |
| $\langle 0,0,1,0,0,1,1,0,0,1,1,0,1\rangle$ | 7 |
| $\langle 1,0,0,1,0,0,1,1,0,0,1,1,0\rangle$ | 5 |
| $\langle 0,1,0,0,1,0,0,1,1,0,0,1,1\rangle$ | 7 |
| $\langle 1,0,1,0,0,1,0,0,1,1,0,0,1\rangle$ | 9 |
| $\langle 1,1,0,1,0,0,1,0,0,1,1,0,0\rangle$ | 3 |
| $\langle 0,1,1,0,1,0,0,1,0,0,1,1,0\rangle$ | 5 |

Table 3.1:

Clearly, no bitstring in $\operatorname{cl}\left(A^{c}\right)$ differs from $A$ in exactly one bit, and therefore, $A$ is not an l-vertex. A similar example can be constructed to show that there exist palindromes which are not l-vertices. However, not all is lost, as can be seen from the next Lemma. This Lemma will be used in the proof of Theorem 3.4.11, which is another useful characterization of $l$-vertices.

Lemma 3.4.10 If some vertex in $c l(A)$ is an anti-palindrome and some other vertex in $\operatorname{cl}(A)$ is a palindrome, then $A$ is an l-vertex.

Proof: Without loss of generality we can assume that $A=\left\langle a_{1}, a_{2}, \ldots a_{n}\right\rangle$ is a palindrome and that $s h^{j}(A)$ in an anti-palindrome. So $a_{i}=a_{n+1-i}$ for all $i$ and $a_{i-j}=$ $1-a_{n+1 \cdots i+j}$ for $i-j \not \equiv_{n} k+1$. We have:

$$
\begin{aligned}
a_{i-j} & =1-a_{n+1-i+j} \text { for } i-j \not \equiv_{n} k+1, \\
a_{i} & =1-a_{n+1-i+2 j} \text { for } i \not \equiv_{n} k+1+j, \\
a_{n-i} & =1-a_{n+1-i+2 j} \text { for } i \not \equiv_{n} k+1+j .
\end{aligned}
$$

Now, if we let $x=n-i$ and $2 j+1=-y$, we see that $a_{x}=a_{x-y}$ for $x \not \equiv_{n} k-j$. Hence A is an $l$-vertex, by Theorem 3.4.4.

THEOREM 3.4.11 $A$ is an $l$-vertex if and only if some vertex in $c l(A)$ is an antipalindrome and some other vertex in $\operatorname{cl}(A)$ is a palindrome.

Proof: Follows from Lemmas 3.4.7, 3.4.8 and 3.4.10.

The following is equivalent to a Lemma that can be found in [4].
Lemma 3.4.12 The number of l-classes in $M_{k}$ is equal to $\phi(n) / 2$.
Proof: Let $A=(q, j, x)=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$ and let $B=(q+j,-j, x)=\left\langle b_{1}, b_{2}, \ldots, b_{n}\right\rangle$.
We claim that $A=B$. We know

$$
\begin{aligned}
& a_{q}=1-a_{q-j}=a_{q-2 j}=1-a_{q-3 j}=\ldots=a_{q+j}, \\
& b_{q+j}=1-a_{q+2 j}=b_{q+3 j}=1-b_{q+4 j}=\ldots=b_{q} .
\end{aligned}
$$

Now, if we write the string of equalities backwards, we get $b_{q}=1-b_{q-j}=b_{q-2 j}=$ $1-b_{q-3 j}=\ldots=b_{q+j}$. Notice that $b_{q}=b_{q+j}=x=a_{q}$. Hence the two strings of equalities:

$$
\begin{gathered}
a_{q}=1-a_{q-j}=a_{q-2 j}=1-a_{q-3 j}=\ldots=a_{q+j} \\
b_{q}=1-b_{q-j}=b_{q-2 j}=1-b_{q-3 j}=\ldots=b_{q+j}
\end{gathered}
$$

are identical. Thus $A=(q, j, x)=(q+j,-j, x)=B$. This also implies that the set of triples $\{(q, j, a): q \in[n]\}$ represents the same class as the set of triples $\{(q, n-j, a): q \in[n]\}$. Note that if $a_{i}=1$, then the triple $\left(i, j, a_{i}\right)$ represents a vertex in $Y$, and if $a=0$, then it represents a vertex in $X$.

Claim: There are exactly two triples that represent each $l$-vertex. By proving the pervious claim we showed that each representation $(q, j, x)$ can be paired with $(q+$ $j,-j, x)$ which represents the same vertex. If we assume that some $l$-vertex $B=$ $\left\langle b_{1}, b_{2}, \ldots, b_{n}\right\rangle$ does not have exactly two triples that represent it, then, without loss of generality, $B$ can be represented by any of the four triples $(1, j, 1),(1+j, n-j, 1)$ and $(q, c, 1),(q+c, n-c, 1)$ with $c \not \equiv_{n} j, n-j$. This means that:

$$
\left\{\begin{array} { l } 
{ b _ { 1 } = 1 } \\
{ b _ { i } = 1 - b _ { i - j } \quad \text { for } i \neq 1 - j }
\end{array} \text { and } \left\{\begin{array}{l}
b_{q}=1 \\
b_{i}=1-b_{i-c} \quad \text { for } i \neq q-c
\end{array}\right.\right.
$$

Hence for $i \neq q-c, 1-j$, we have $b_{i}=1-b_{i-c}=1-b_{i-j}$, which implies that $b_{i-c}=b_{i-j}$ for $i \neq q-c, 1-j$. This then implies that $b_{i}=b_{i+c-j}$ for $i \neq q, 1+c-j$. Consider the sequence $b_{1}, b_{1+c-j}, b_{1+2(c-j)}, \ldots, b_{1+(n-1)(c-j)}$. Note that because $b_{i}=b_{i+j-c}$ for $i \neq q, 1+c-j$, the sequence will be $1,1, \ldots 1,0,0, \ldots 0$, where the first zero will occur at
$b_{q+(j-c)}$. We know that B has $k+1$ 1's and $k 0$ 's. So it must be that our sequence $1,1, \ldots 1,0,0, \ldots 0$ has $k+1$ ones followed by $k$ zeros. This gives us $1+(k+1)(j-c) \equiv$ $q+j-c(\bmod n)$. Similarly, if we instead look at $b_{q}, b_{q+c-j}, b_{q+2(c-j)}, \ldots, b_{q+(n-1)(c-j)}$ we get that $q+k(j-c) \equiv 1(\bmod n)$. Combining these two results we obtain:

$$
\begin{aligned}
q+1+(n)(j-c) & \equiv_{n} 1+j-c+q \\
q+1 & \equiv_{n} 1+j-c+q \\
0 & \equiv_{n} j-c .
\end{aligned}
$$

Now, $0 \equiv_{n} j-c$ implies $j \equiv_{n} c$, which is a contradiction.
Each $l$-class in $Y$ can be represented as either the set $\{(q, j, 1): q \in[n]\}$ or $\{(q, n-$ $j, 1): q \in[n]\}$ for some $j$ and each $l$-class in $X$ can be represented as either the set $\{(q, j, 0): q \in[n]\}$ or $\{(q, n-j, 0): q \in[n]\}$ for some $j$. Hence the number of $l$-classes is equal to twice the number of unordered pairs $(j, n-j)$ with $\operatorname{gcd}(j, n)=1$ and $j<n-j$. This is exactly $\phi(n) / 2$.

THEOREM 3.4.13 If $A$ is adjacent to two distinct vertices from the same class, then that class must be the class containing $A^{c}$.

Proof: Suppose that $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$ is adjacent to $B$ and $s h^{j}(B)$ with $1 \leq j \leq k$. From the proof of Lemma 3.4.2 we know that $B=\left\langle b_{1}, b_{2}, \ldots, b_{n}\right\rangle$ is adjacent to $A$ and $s h^{j}(A)$. Without loss of generality, let B differ from $A$ in the $1^{\text {st }}$ bit and from $s h^{j}(A)$ in the $t^{\text {th }}$ bit. Since B differs from $A$ only in the $1^{\text {st }}$ bit, we have

$$
\left\{\begin{array}{l}
a_{i}=b_{i} \quad \text { for } i \neq 1 \\
a_{1}=1-b_{1}
\end{array}\right.
$$

and because $B$ differs from $s h^{j}(A)$ only in the $t^{t h}$ bit, we have
$\left\{\begin{array}{l}a_{i+j}=b_{i} \quad \text { for } i \neq t \\ a_{t+j}=1-b_{t} .\end{array}\right.$
Combining these results we get that $a_{i}=a_{i+j}$ for $i \neq 1, t$. Without loss of generality, we can assume that A has $k+11$ 's. Using a similar argument to the one we used in Lemma 3.4.12, we obtain $1+(k+1)(j) \equiv t(\bmod n)$ and
$a_{1+j}=a_{1+2 j}=\ldots a_{1+(k+1) j}=1-a_{1+(k+2) j}=1-a_{1+(k+3) j}=\ldots=1-a_{1}$.
Note that this implies $a_{i}=1-a_{i+(k+1) j}$ for all $i \neq 1$. Hence A can be represented by the triple $(1,(k+1) j, 1)$. Now, repeating this argument for $B$, we see that $B$ can be represented by $(t,(k) j, 0)$. Since $k j \equiv_{n} n-(k+1) j$, we establish that B is in the complement class of A.

Lemma 3.4.14 The vertex $(s, j, f)$ is adjacent to the vertices $(s-j, j, 1-f)$ and $(s+j, j, 1-f)$.

Proof: Let $A=\left\langle a_{1}, a_{2}, \ldots a_{n}\right\rangle=(s, j, f)$ and $B=\left\langle b_{1}, b_{2}, \ldots b_{n}\right\rangle=(s-j, j, 1-f)$. We know that $f=a_{s}=1-a_{s-j}=a_{s-2 j}=\ldots=a_{s-(n-1) j}$ and that $1-f=b_{s-j}=$ $1-b_{s-2 j}=b_{s-3 j}=\ldots=b_{s-n j}$. Note that $s-n j \equiv_{n} s$. From 1-f= $b_{s-j}=1-b_{s-2 j}=$ $b_{s-3 j}=\ldots=b_{s-n j}$ we get that $f=1-b_{s-j}=b_{s-2 j}=1-b_{s-3 j}=\ldots=1-b_{s}$. Thus $A$ differs from $B$ in only the $s^{\text {th }}$ bit. Showing that $(s, j, f)$ is adjacent to $(s+j, j, 1-f)$ is similar.

If we wish to build a string which represents the $l$-vertex A with triple $\left(s, j, a_{s}\right)$, then the 1's and 0's get added in a certain order. For example, if we wanted to build the string representing the $l$-vertex $A \in V\left(M_{5}\right)$ with triple $(1,3,1)$, then we would start with a blank bitstring of length $5:(*, *, *, *, *)$. We know that $a_{1}=1$, so that would be our first bit added, giving us $(1, *, *, *, *)$. Next we know that $a_{1}=1-a_{1-3}=1-a_{3}$, so $a_{3}$ is our second bit added, giving us $(1, *, 0, *, *)$. Now, $a_{3}=1-a_{3-3}=1-a_{5}$, so $a_{5}$, is our third bit added, giving us $(1, *, 0, *, 1)$. We continue in this manner until we have the entire string:
$(*, *, *, *, *) \rightarrow(1, *, *, *, *) \rightarrow(1, *, 0, *, *) \rightarrow(1, *, 0, *, 1) \rightarrow(1,0,0, *, 1) \rightarrow(1,0,0,1,1)$
The above gives us the order in which the bits of an $l$-vertex get added to the string. This means that we can represent a non-l-vertex which is adjacent to an $l$-vertex by just indicating in which bit it differs from its adjacent $l$-vertex. We will use the notation $\left(s, j, a_{s},[i]\right)$, to refer to such a vertex. Here $[i]$ indicates that the $i^{\text {th }}$ bit added is the bit in which this vertex differs from the $l$-vertex. Similarly, we can introduce notation for any vertex in $M_{k}$ by listing the positions in which it differs from some $l$-vertex. Define the notation to be $\left(s, j, a_{s},\left[r_{1}, r_{2}, . ., r_{c}\right]\right)$, where $\left[r_{1}, r_{2}, . ., r_{c}\right]$ is a list of bits. For simplicity we will stipulate that $r_{i} \neq r_{j}$ when $i \neq j$, because any two strings either differ or do not differ in any given bit.

Observation 3.4.15 In say, $\left(s, j, f,\left[u_{1}, u_{2}, \ldots u_{r}\right]\right)$, note that $\left[u_{1}, u_{2}, \ldots u_{r}\right]$ is just a set of instructions that tell us which bits of $(s, j, f)$ to change in order to get $\left(s, j, f,\left[u_{1}\right.\right.$, $\left.u_{2}, \ldots u_{r}\right]$ ). If we change $s$ or $j$, the $i^{\text {th }}$ bit we add will still have the same value. This
means that $\left(s_{1}, j_{1}, f_{1},\left[u_{1}, u_{2}, \ldots u_{c}\right]\right)$ is adjacent to $\left(s_{1}, j_{1}, f_{1},\left[v_{1}, v_{2}, \ldots v_{r}\right]\right)$ if and only if $\left(s_{2}, j_{2}, f_{2},\left[u_{1}, u_{2}, \ldots u_{c}\right]\right)$ is adjacent to $\left(s_{2}, j_{2}, f_{2},\left[v_{1}, v_{2}, \ldots v_{s}\right]\right)$. Also, $\left(s_{1}, j_{1}, f_{1},\left[u_{1}, u_{2}, \ldots\right.\right.$ $\left.\left.u_{c}\right]\right)=\left(s_{1}, j_{1}, f_{1},\left[v_{1}, v_{2}, \ldots v_{r}\right]\right)$ if and only if $\left(s_{2}, j_{2}, f_{2},\left[u_{1}, u_{2}, \ldots u_{c}\right]\right)=\left(s_{2}, j_{2}, f_{2},\left[v_{1}, v_{2}, \ldots\right.\right.$ $\left.v_{r}\right]$ ).

Lemma 3.4.16 Let $A$ be the vertex represented by the triple $(s, j, f)$. Then $(s, j, f,[1])$ and $(s, j, f,[n])$ are in the same class as $A^{c}$.

Proof: Let $A=(s, j, f)=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$. Without loss of generality we can assume that $f=1$. This gives us

$$
\begin{cases}a_{s}=1 & \\ a_{s-(2 i-1) j}=1-a_{s}=0, & 1 \leq i \leq k \\ a_{s-2 i j}=a_{s}=1, & 1 \leq i \leq k\end{cases}
$$

This means that for $(s, j, f,[1])$ we have

$$
\begin{cases}a_{s}=0 \\ a_{s-(2 i-1) j}=0, & 1 \leq i \leq k \\ a_{s-2 i j}=1, & 1 \leq i \leq k\end{cases}
$$

which is equivalent to:

$$
\begin{cases}a_{(s+j)}=0 & \\ a_{(s+j)-(2 i-1) j}=1=1-a_{(s+j)}, & 1 \leq i \leq k \\ a_{(s+j)-2 i j}=0=a_{(s+j)}, & 1 \leq i \leq k\end{cases}
$$

meaning that $(s, j, f,[1])$ represents the same vertex as $(s+j, j, 1-f)$, which gives us the desired result. The case $(s, j, f,[n])$ is done similarly.

Take some vertex $A=\left\langle a_{1}, a_{2}, \ldots a_{n}\right\rangle$ and some $x \in\{2,3, . . n-1\}$ which is relatively prime to $n$ and form the vertex $B=\left\langle b_{1}, b_{2}, \ldots b_{n}\right\rangle$ by defining $b_{i}=a_{x^{-1}}$, or equivalently $a_{i}=b_{x i}$, for $1 \leq i \leq n$. (Since $x$ is relatively prime to $n, x$ is invertible modulo $n$ ). We call this action multiplying and we write $\operatorname{mult}_{x}(A)=B$.

Observation 3.4.17 Let $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle=(i, j, f)$ with $n$ prime. We know that $a_{i}=1-a_{i-j}=a_{i-2 j}=1-a_{i-3 j}=\ldots=a_{i-(n-1) j}$. So when we take $\operatorname{mult}_{x}(A)=A^{\prime}=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, \ldots, a_{n}^{\prime}\right\rangle$, we have that $a_{x i}^{\prime}=1-a_{x i-x j}^{\prime}=a_{x i-2 x j}^{\prime}=1-$ $a_{x i-3 x j}^{\prime}=\ldots=a_{i x-(n-1) x j}^{\prime}$. This gives us $A^{\prime}=(i x, j x, f)$. Similarly, we can establish that $\operatorname{mult}_{x}\left(\left(i, j, f\left[r_{1}, r_{2}, . . r_{s}\right]\right)\right)=\left(x i, x j, f\left[r_{1}, r_{2}, . . r_{s}\right]\right)$. So if mult ${ }_{x}(A)=B$, then for every $A^{\prime} \in \operatorname{cl}(A), \operatorname{mult}_{x}\left(A^{\prime}\right) \in \operatorname{cl}(B)$.

Note 3.4.18 The vertices $A$ with the property that $\operatorname{sh}^{j}\left(\operatorname{mult}_{x}(A)\right)=A$ for some $j$ and some $x \not \equiv_{n}-1,0,1$ will have a special role later. We will refer to these as central vertices. Note that a consequence of Observation 3.4.17 is that if $A$ is a central vertex, then so is every other vertex in its class as well as all the vertices in its inverse class.

THEOREM 3.4.19 If $A$ is a central vertex with $\operatorname{sh}^{j}\left(\operatorname{mult}_{x}(A)\right)=A$, then for every $y \in\langle x\rangle$ there exists $j$ such that $\operatorname{sh}^{j}\left(\right.$ mult $\left._{y}(A)\right)=A$.

Proof: By the remark preceding this theorem, if $A$ is a central vertex, then so is every other vertex in $c l(A)$. We know that mult $_{x}(A) \in c l(A)$. Combining these results, we get that mult $_{x}\left(\right.$ mult $\left._{x}(A)\right)=$ mult $_{x^{2}}(A) \in c l(A)$. We can repeat this argument to get that mult $_{x^{m}}(A) \in c l(A)$ for all $m$. Every element in $\langle x\rangle$ can be written as $x^{m}$ for
some $m$. Hence for every $y \in\langle x\rangle$ we have $\operatorname{mult}_{y}(A) \in \operatorname{cl}(A)$. Therefore, for every $y \in\langle x\rangle$ there exists a $j$ such that $s h^{j}\left(\operatorname{mult}_{y}(A)\right)=A$.

Lemma 3.4.20 For any vertex $A \in V\left(M_{k}\right)$, mult $t_{n-1}(A)=f l(A)$.
Proof: Let $A^{\prime}=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, \ldots a_{n}^{\prime}\right\rangle=$ mult $_{n-1}(A)$. Observe that $j(n-1) \equiv j n-j \equiv$ $n-j(\bmod n)$. This means that $a_{n-i}^{\prime}=a_{i}$ for all $i$, and therefore $\operatorname{mult}_{n-1}(A)=f l(A)$, as desired.

Observation 3.4.21 From observation 3.4.17, we know that mult ${ }_{x}(i, j, f)=(i x, j x$ ,f). We wish to show that if $A$ is central and $x \neq n-1$, then $\operatorname{mult}_{x}(A) \notin \operatorname{cl}(A)$. If we assume toward a contradiction that $j x \equiv j(\bmod n)$, we get that $x \equiv_{n} 1$, which is a contradiction because, $x \in\{2,3, . . n-1\}$. So $j x \not \equiv j(\bmod n)$. If we assume that $j x \equiv n-j(\bmod n)$, we get $x \equiv_{n} n-1$. In Lemma 3.4.20 we established that mult $_{n-1}\left(\left(i, j, f\left[r_{1}, r_{2}, . . r_{s}\right]\right)\right)=f l\left(\left(i, j, f\left[r_{1}, r_{2}, . . r_{s}\right]\right)\right)$. Therefore, there is only one multiplication of a vertex which sends that vertex to a vertex in the same class and that is the one which flips it. Another consequence of this is that if $k>2$, then no vertex is both a central vertex and an l-vertex.

Lemma 3.4.22 If $A$ is a central vertex with $\operatorname{sh}^{j}\left(\operatorname{mult}_{x}(A)\right)=A$, then

$$
s h^{j+1-x}\left(m u l t_{x}(\operatorname{sh}(A))\right)=\operatorname{sh}(A) .
$$

Proof: Let $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$ be a central vertex with $\operatorname{sh}^{j}\left(\operatorname{mult}_{x}(A)\right)=A$ and let
$A^{\prime}=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, \ldots, a_{n}^{\prime}\right\rangle$ be the vertex $\operatorname{sh}(A)$. We know that $a_{x i-j}=a_{i}$ for all $i$. Also, $a_{i}^{\prime}=a_{i-1}$ for all $i$. So $a_{x i-x-j+1}^{\prime}=a_{x(i-1)-j}=a_{i-1}=a_{i}^{\prime}$. This gives $a_{x i+(j-x+1)}^{\prime}=a_{i}^{\prime}$ for all $i$, which implies that $s h^{j-x+1}\left(\operatorname{mult}_{x}\left(A^{\prime}\right)\right)=A^{\prime}$, as desired.

Corollary 3.4.23 If $n$ is prime and $A$ is a central vertex in $M_{k}$ with $\operatorname{sh}^{j}\left(\operatorname{mult}_{x}(A)\right)=$ $A$, then for every $r$ there is some vertex $A^{\prime} \in \operatorname{cl}(A)$ with $\operatorname{sh}^{r}\left(\operatorname{mult}_{x}\left(A^{\prime}\right)\right)=A^{\prime}$.

Proof: Let $n$ be a prime and $A$ a central vertex in $M_{k}$ with $s h^{j}\left(\operatorname{mult}_{x}(A)\right)=A$ and let $r$ be fixed. By repeatedly applying Lemma 3.4.22, we get $s h^{j+y(1-x)}\left(\operatorname{mult}_{x}\left(s h^{y}(A)\right)\right)=$ $\operatorname{sh}^{y}(A)$ for all $y$. Now, $j+y(1-x) \equiv_{n} r$ for some $y$. Solving for $y$, we get $y \equiv_{n}(r-j)(1-x)^{-1}$. If follows that $s h^{(r-j)(1-x)^{-1}}(A)$ has the desired property.

Lemma 3.4.24 If $n$ is prime and $A$ is adjacent to $A^{\prime}$ then $\operatorname{mult}_{x}(A)$ is adjacent to $\operatorname{mult}_{x}\left(A^{\prime}\right)$ for every $x \in\{2,3, . ., n-1\}$.

Proof: Let $A=\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$ and $A^{\prime}=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, \ldots, a_{n}^{\prime}\right\rangle$ be adjacent. Then $A$ differs from $A^{\prime}$ in one bit, say $a_{i}$. Let $x \in\{2,3, . ., n-1\}$. Since $n$ is prime and $x \in$ $\{2,3, . ., n-1\}$ we know that $x$ has a multiplicative inverse modulo $n$, namely $x^{-1}$. Hence $\operatorname{mult}_{x}(A)=\left\langle a_{1 x^{-1}}, a_{2 x^{-1}}, \ldots, a_{n x^{-1}}\right\rangle$ and $\operatorname{mult}_{x}(A)=\left\langle a_{1 x^{-1}}^{\prime}, a_{2 x^{-1}}^{\prime}, \ldots, a_{n x^{-1}}^{\prime}\right\rangle$. However, we know that $a_{j}=a_{j}^{\prime}$ for all $j \not \mathcal{F}_{n} i$, so that $\operatorname{mult}_{x}(A)$ and $m u l t_{x}\left(A^{\prime}\right)$ differ in only the $i x^{t h}$ bit and are therefore adjacent.

Observation 3.4.25 Lemma 3.4.24 showed us that if $n$ is prime and $A$ is adjacent to $A^{\prime}$ then $\operatorname{mult}_{x}(A)$ is adjacent to mult $_{x}\left(A^{\prime}\right)$ for every $x \in\{2,3, \ldots, n-1\}$. Hence the
mapping from $V\left(M_{k}\right)$ to itself, given by mult $_{x}$, preserves edges and it is therefore a homomorphism. Since the mapping is invertible (mult $t_{x^{-1}}$ is its inverse) this mapping is in fact an automorphism.

### 3.5 Properties of $R_{k}$

Déjter [4] observed a useful reduction of the the problem. The problem of finding a Hamiltonian cycle in $M_{k}$ can be reduced to the problem of finding a Hamiltonian path in a smaller graph, $R_{k} . R_{k}$ is obtained in the following way:

First we obtain the reduced graph $N_{k}$. Let each vertex in $N_{k}$ represent a distinct class of $M_{k}$ and let two vertices in $N_{k}$ be adjacent if their corresponding classes in $M_{k}$ contain vertices which are adjacent. Two vertices in $N_{k}$ will be called inverses of each other if they represent classes which are inverses of one another. The reduced graph $R_{k}$ is obtained by letting each vertex represent a pair of vertices which are inverses of one another in $N_{k}$. Two vertices in $R_{k}$ are adjacent if they represent adjacent vertices in $N_{k}$, and if a vertex in $N_{k}$ is adjacent to its inverse, the corresponding vertex will have a loop in $R_{k}$. Now, because of how we construct $R_{k}$, each vertex in it represents $2 n$ vertices in two classes which are inverses of one another in $M_{k}$. We will refer to these $2 n$ vertices in $M_{k}$ as the parents of the vertex they correspond to in $R_{k}$. Notice that vertices in $R_{k}$ which have loop edges are exactly the vertices whose parents are $l$-vertices ( $l$ for loop). Déjter showed that a Hamiltonian path in $R_{k}$ with loop vertices as endpoints can be used to construct a Hamiltonian cycle in $M_{k}$.

Observation 3.5.1 The representations for the vertices of $M_{k}$ that we used earlier can be carried over to $R_{k}$; we can even further simplify the notation: For example, in $\left(i, j, f\left[u_{1}, u_{2}, \ldots u_{s}\right]\right)$, the $i$ and $f$ are redundant, so we can represent the same vertex with $\left(j,\left[u_{1}, u_{2}, \ldots u_{s}\right]\right)$. Note that the l-vertex $(i, j, f)$ will now be represented as $(j)$.

The following is a special case of a Theorem proven in [4].

THEOREM 3.5.2 If there exist a Hamiltonian path from (1) to (k) in $R_{k}$, then there exists a Hamiltonian cycle in $M_{k}$.

Proof: Using the relationship between $R_{k}$ and $M_{k}$ we can turn a Hamiltonian path in $R_{k}$ into $2 n$ disjoint paths which together include all of the vertices of $M_{k}$.


Figure 3.3:

Now, all we need to do is connect the paths into a cycle. We add the $n$ edges of the from $(j, 1,1)(j-1,1,0)$ (we showed in Lemma 3.4.14 that these edges exist). Now we have $n$ disjoint paths which together include all of the vertices of $M_{k}$.


Figure 3.4:

We know we have $n$ edges of the from $(j, k, 1)(j+k, k, 0)$. We claim that adding these completes the cycle.


Figure 3.5:

It suffices to show that we can get from any vertex of the form $(j, k, 1)$ to any other vertex of that form. We know that $(j, k, 1)$ is connected through a path to $(j-1, k, 0)$
which is connected to $(j-k-1, k, 1)=(j-(k+1), k, 1)$. Repeating this we can get from $(j, k, 1)$ to $(j-x(k+1), k, 1)$ for any $x$. Since $\operatorname{gcd}(k+1, n)=1$, the integers $j, j-(k+1), j-2(k+1), \ldots, j-(n-1)(k+1)$ form a complete system of residues modulo $n$. Thus we can get from any vertex of the form $(j, k, 1)$ to any other vertex of that form.

We now illustrate the construction in Theorem 3.5 .2 by example:

Let us start with $R_{3}$ pictured below. We construct a Hamiltonian path between two loop vertices $(1)=1010101,(3)=1111000$. This path can be explicitly written as: 1010101, $1010001,1011001,1011000,1111000$ and is represented below as a colored path.


Figure 3.6:

This path in $R_{3}$ represents $2(7)=14$ disjoint paths in in $M_{3}$. The edges are color matched so one can see which edges in $R_{3}$ correspond to edges in $M_{3}$. Now add all edges of the form $(j, 1,1)(j-1,1,0)$ as well as all edges of the form $(j, k, 1)(j+k, k, 0)$. This, along with the paths we already have, gives us a Hamiltonian cycle in $M_{k}$.


Figure 3.7:

If the reader wishes to follow the path directly, the following diagram illustrates it with the edges color matched as in the previous diagrams.


Figure 3.8:

Corollary 3.5.3 Let $n$ be prime. If there exists a Hamiltonian path in $R_{k}$ from (a) to (b) and $g c d(a, b)=1=\operatorname{gcd}(a+b, n)$, then there exists a Hamiltonian cycle in $M_{k}$. The proof of this result is essentially identical to the proof of Theorem 3.5.2 with $a$ replacing 1 and $b$ replacing $k$.

Lemma 3.5.4 $R_{k}$ contains $\frac{(n-1)!}{(k)!(k+1)!}$ vertices.
Proof: Each class has $n$ vertices in it and each class has a corresponding unique inverse class, so that each vertex in $R_{k}$ represents exactly $2 n$ vertices in $M_{k}$. Since there are $2\binom{n}{k}$ vertices in $M_{k}$, then the number of vertices in $R_{k}$ is $2 \times\binom{ n}{k} /(2 n)=$ $\frac{(n-1)!}{(k)!(k+1)!}$, as claimed.

Dejter noted that the number of vertices in $R_{k}$ is always a Catalan number. We will explore the relationship between $R_{k}$ and the Catalan numbers more thoroughly in the next section.

Observation 3.5.5 The number of loop vertices in $R_{k}$ is exactly $\phi(n) / 2$. This follows directly from Lemma 3.4.12.

Observation 3.5.6 If $n$ is prime then, $R_{k}$ has $k+1$ vertices of degree $k-1$ and $\frac{(n-1)!}{(k)!(k+1)!}-(k+1)$ vertices of degree $k+1$.

When reducing $M_{k}$ to $R_{k}$, the two most natural automorphisms, shifting and flipping the string are used to condense $M_{k}$ into $R_{k}$. A natural question is whether there are any more automorphisms present in $R_{k}$, and if so, can we use them to reduce the $R_{k}$ further into an even smaller graph? The answer to this question is positive in some special cases. To illustrate these, we need to develop some further machinery.

We want to extend the definition of multiplication to $R_{k}$ : For $A, B \in R_{k}$ we say that $\operatorname{mult}_{x}(A)=B$ if for every parent of $A$, say, $A^{\prime} \in M_{k}, m u l t_{x}\left(A^{\prime}\right)$ is a parent of $B$. Using this definition we can now extend the definition of central vertices to
include vertices in $R_{k}$. That is: $A \in R_{k}$ is a central vertex if $\operatorname{mult}_{x}(A)=A$ for some $x \not \equiv_{n}-1,0,1$. We have already observed that if $A \in M_{k}$ is a central vertex, then so is every other vertex in both its class and its inverse class. This means that central vertices in $R_{k}$ have parents which are central vertices in $M_{k}$.

Lemma 3.5.7 Let $n=2 k+1$ be prime. If a vertex $A \in R_{k}$ can be represented as both $\left(i,\left[u_{1}, u_{2}, . . u_{s}\right]\right)$ and $\left(j,\left[u_{1}, u_{2}, . . u_{s}\right]\right)$ with $i \not \equiv_{n} j, n-j$, then $A$ can be represented as $\left(r,\left[u_{1}, u_{2}, . . u_{s}\right]\right)$ for any $r \in\left\langle i j^{-1}\right\rangle$.

Proof: This is a consequence of Note 3.4.18 and Theorem 3.4.19: If $A \in M_{k}$ is central, then for every $r \in\left\langle i j^{-1}\right\rangle$ there exists an $m$ such that $s h^{m}\left(\operatorname{mult}_{r}(A)\right)=A$.

THEOREM 3.5.8 Let $k>3$. If both $k$ and $n=2 k+1$ are prime, then there is at least one central vertex in $M_{k}$.

Proof: The set $\mathbb{Z}_{n} \backslash\{0\}$ is a group under multiplication. This groups has order $2 k$. By the first Sylow theorem, we know that this group has some subgroup $\mathbb{G}$ of order $k$. Let the the elements of $\mathbb{G}$ be $\left\{x_{1}, x_{2}, \ldots x_{k}\right\}$. Now let $A=\left\langle a_{1}, a_{2}, . ., a_{n}\right\rangle$ with $a_{i}=1$ if $i \in\left\{x_{1}, x_{2}, \ldots x_{k}\right\}$ and $a_{i}=0$ otherwise (note that $A$ has $k 1$ 's and $k+10$ 's as desired). We claim that $A$ is central. Consider $A^{\prime}=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, . ., a_{n}^{\prime}\right\rangle=\operatorname{mult}_{x_{i}}(A)$. The cosets of the subgroup $\mathbb{G}$ are either disjoint or identical. So $j \mathbb{G}=\mathbb{G}$ if $j \in \mathbb{G}$ and $j \mathbb{G} \cap \mathbb{G}=\emptyset$ otherwise. Hence $j x_{i}^{-1} \in\left\{x_{1}, x_{2}, \ldots x_{k}\right\}$ if and only if $j \in\left\{x_{1}, x_{2}, \ldots x_{k}\right\}$, since $x_{i}^{-1} \in \mathbb{G}$. Thus, $a_{j}^{\prime}=a_{j x_{i}^{-1}}=1$ if and only if $j \in\left\{x_{1}, x_{2}, \ldots x_{k}\right\}$, and $a_{j}^{\prime}=0$ otherwise. So mult $_{x_{i}}(A)=A$. In fact, mult $_{x_{i}}(A)=A$ regardless of which $x_{i}$ we choose. This establishes that $A$ is a central vertex, because at least one element of
$\mathbb{G}$ is not $-1,0,1$.

THEOREM 3.5.9 If $k$ and $n=2 k+1$ are both prime, then there are exactly two central vertices in $R_{k}$ and these two vertices are adjacent to each other.

Proof: Let $A$ be a central vertex in $R_{k}$ and $\mathbb{G}$ the group from the previous theorem.
From Corollary 3.4.23 and Theorem 3.4.19, we know that for any pair $b \in \mathbb{G}$ and $c \in$ $\{0,1,2,3, \ldots n-1\}$ there exist a parent of $A$, say $A^{\prime} \in M_{k}$, such that $\operatorname{sh}^{c}\left(m u l t_{b}\left(A^{\prime}\right)\right)=$ $A^{\prime}=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, . . a_{n}^{\prime}\right\rangle$. This means that for any $i, a_{b i-c}^{\prime}=a_{i}^{\prime}$. Let $f(x)=b x-c$. We know that for any fixed $i, a_{i}^{\prime}=a_{f(i)}^{\prime}=a_{f^{2}(i)}^{\prime}=a_{f^{3}(i)}^{\prime}=\ldots$. Since $f^{r}(i)$ can take on a finite number of values, there exist some $0 \leq r_{1}<r_{2}<n$ such that $f^{r_{1}}(i)=f^{r_{2}}(i)$. Without loss of generality we can assume that $r_{2}$ is minimum. The function $f(x)$ is invertible: $f^{-1}(x)=b^{-1} x+b^{-1} c$. So if $r_{1} \neq 0$ :

$$
\begin{aligned}
& f^{r_{1}}(i) \equiv_{n} f^{r_{2}}(i) \\
& f^{-r_{1}}\left(f^{r_{1}}(i)\right) \equiv_{n} f^{-r_{1}}\left(f^{r_{2}}(i)\right) \\
& i \equiv_{n} f^{r_{2}-r_{1}}(i)
\end{aligned}
$$

This contradicts the minimality of $r_{2}$ since $r_{2}-r_{1}<r_{2}$. So $r_{1}$ must be zero, meaning that $i \equiv{ }_{n} f^{r_{2}}(i)$. This means $f^{j}(i) \equiv_{n} f^{r_{2}+j}(i)$ for all $j$. If $r_{2}>k$ then we have more that $k+1$ bits all equal to one another. This is a contradiction, because $A^{\prime}$ has either $k+11$ 's and $k 0$ 's or $k+10$ 's and $k 1$ 's. So $0<r_{2} \leq k$. If we assume that $r_{2}=1$, then we can solve for $i: i \equiv_{n} b i-c \rightarrow i \equiv_{n} c(b-1)^{-1}$. Since we can solve for it, there
is exactly one value of $i$ for which $r_{2}$ is 1 . If we assume that $r_{2}=2$ we get

$$
\begin{gathered}
i \equiv_{n} f^{2}(i) \\
i \equiv_{n} b^{2} i-c b-c \\
c(b+1) \equiv_{n}\left(b^{2}-1\right) i \\
\left(b^{2}-1\right)^{-1} c(b+1) \equiv_{n} i \\
\left(b^{2}-1\right)^{-1} c(b+1) \equiv_{n} i \\
(b+1)^{-1}(b-1)^{-1} c(b+1) \equiv_{n} i \\
(b-1)^{-1} c \equiv_{n} i
\end{gathered}
$$

So in this case we notice that $i \equiv_{n} c(b-1)^{-1}$ happens to be the same $i$ we got when solving for $i$ when $r_{2}=1$. Thus the minimum $r$ is never 2 for any $i$. Claim: $f^{m}(x)=b^{m} x-c\left(b^{m-1}+b^{m-2}+. .+b+1\right)$. We proceed to show this by induction. The base case, $f(x)=f^{1}(x)=b x+c$, is given. Now, assume that the claim holds for all $m<M-1$ :

$$
\begin{aligned}
& f^{M}(x)=f\left(f^{M}(x)\right) \\
& f^{M}(x)=f\left(b^{M-1} x-c\left(b^{M-2}+b^{M-3}+. .+b+1\right)\right. \\
& f^{M}(x)=b\left(b^{M-1} x-c\left(b^{M-2}+b^{M-3}+. .+b+1\right)+c\right. \\
& f^{M}(x)=\left(b^{M} x-c\left(b^{M-1}+b^{M-1}+. .+b\right)\right)+c \\
& f^{M}(x)=b^{M} x-c\left(b^{M-1}+b^{M-1}+. .+b+1\right)
\end{aligned}
$$

Thus $f^{m}(x)=b^{m} x-c\left(b^{m-1}+b^{m-2}+. .+b+1\right)$. Note that $b^{2 k} \equiv_{n} 1$ by Fermat's Little Theorem. So we have:

$$
\begin{aligned}
f^{2 k}(x) & =b^{2 k} x-c\left(b^{2 k-1}+b^{2 k-2}+. .+b+1\right) \\
& =b^{2 k} x-c\left(b^{2 k}-1\right)(b-1)^{-1} \\
& \equiv_{n}(1) x-c(1-1)(b-1)^{-1} \\
& =x
\end{aligned}
$$

Thus $f^{2 k}(x)=x$. This means that $r_{2} \mid 2 k$. We showed that $r_{2} \neq 2$ and that $r_{2} \leq k$ so $r_{2}$ can only be 1 or $k$ since we assumed that $k$ was prime. So $r_{2}$ is 1 when $i=c(b-1)^{-1}$, and $k$, otherwise. So we have 3 sets of bits of sizes $1, k$ and $k$ each of which must be either all 1's or all 0 's. Since there are either $k$ or $k+11$ 's, one of the sets of size $k$ has to be all 1's and one of the sets has to be all 0 's. If the bit in set of size one is 1 , then we get one of the central vertices, and if it is 0 , we get a neighboring central vertex. Observation 3.4.21 tells us that $l$-vertices are not central vertices. So these two adjacent central vertices must be distinct parents of two distinct central vertices in $R_{k}$. Thus we have constructed two central vertices in $R_{k}$. By Corollary 3.4.23 and Theorem 3.4.19, we know that had we had chosen any other $b, c$ we would have gotten the same two central vertices. So they must be unique.

### 3.6 The Automorphisms of $R_{k}$

Let $k$ and $n$ both be prime. The set $\mathbb{Z}_{n} \backslash\{0\}$ is a group of order $2 k$ under multiplication. The first Sylow Theorem (Theorem 2.2.7) tells us that $\mathbb{Z}_{n} \backslash\{0\}$ has a subgroup $\mathbb{G}$ of order $k$. We know show that this subgroup is unique.

THEOREM 3.6.1 If $n$ and $k$ are both prime then, $\mathbb{Z}_{n} \backslash\{0\}$ has exactly one subgroup of order $k$.

Proof: We know, by the Third Sylow Theorem (Theorem 2.2.8), that the number of such subgroups is congruent to 1 modulo $k$. So if we assume toward at contradiction that there is more than one such subgroup, then there must be at least $k+1$ of them. At least two of these subgroups must have a common element, $y$, other than 1 , by the pigeonhole principle. However, since $k$ is prime, these subgroups are cyclic and hence both equal to $\langle y\rangle$, a contradiction. So $\mathbb{G}$ is unique.

The mapping from $\phi: R_{k} \mapsto R_{k}$ given by $A \mapsto \operatorname{mult}_{x}(A)$ is an automorphism for all $x \in\{1,2,3 \ldots, n-1\}$. Let $A=\left\langle a_{1}, a_{2}, . ., a_{n}\right\rangle$ with $a_{i}=1$ if and only if $i \in \mathbb{G}$. We showed in the proof of Theorem 3.5.8 that if $x \in \mathbb{G}$, then $A=\operatorname{mult}_{x}(A)$. However, if $x \notin \mathbb{G}+\{0\}$, then $i x \in \mathbb{G}$ if and only if $i \notin \mathbb{G}+\{0\}$. This means that when $x \notin \mathbb{G}+\{0\}$ and $C$ and $C^{\prime}$ are the two central vertices, mult $_{x}(C)=C^{\prime}$ and mult $_{x}\left(C^{\prime}\right)=C$. This happens because $\operatorname{mult}_{x}\left(C^{c}\right)$ differs from $C$ in exactly the $0^{t h}\left(=n^{t h}\right)$ bit. Hence the mapping given by $\operatorname{mult}_{x}(C)$ with $x \in \mathbb{G}$ fixes the central vertices and when $x \notin \mathbb{G}+\{0\}$ it maps them to one another. So we will refer to the $k$ automorphisms that fix the central vertices as central fixing automorphisms, or $C F A$ 's for short.

Lemma 3.6.2 If $n$ and $k>2$ are both prime, then the unique subgroup of $\mathbb{Z}_{n} \backslash\{0\}$ of order $k$ is $\left\langle x^{2}\right\rangle$, where $x$ is an arbitrary element in $\mathbb{Z}_{n} \backslash\{-1,0,1\}$.

Proof: The only element of order 2 in $\mathbb{Z}_{n} \backslash\{0\}$ is -1 . Lagrange's Theorem tells us that the order of an element divides the order of the group. Thus if $x \in \mathbb{Z}_{n} \backslash\{-1,0,1\}$ then, $x$ has either order $2 k$ or $k$. It suffices to show that $x^{2}$ is of order $k$. By Ferrmat's Little Theorem $\left(x^{2}\right)^{k} \equiv_{n} x^{2 k} \equiv_{n} x^{n-1} \equiv_{n} 1$. Hence $\left(x^{2}\right)^{k} \equiv_{n} 1$ and $x^{2}$ has order $k$.

The above Lemma shows us that, we can use any square to generate the subgroup of order $k$. We showed the uniqueness of this subgroup in Theorem 3.6.1. Using the proofs of Theorems 3.5.8 and 3.5.9 we can then use this subgroup to generate all central vertices in $R_{k}$.

Let's explore the automorphisms of $R_{3} . R_{3}$ has 5 vertices. Three of them are loop vertices and the other two are central vertices. We label the three loop vertices (1), (2) and (3) and the two central vertices $C_{1}$ and $C_{2}$. The automorphism represented by the diagram is written underneath it. mult $t_{1}$ is, of course, the identity. Since the unique subgroup of $\mathbb{Z}_{7} \backslash\{0\}$ of order 3 is $\{1,2,4\}$ mult $_{1}$, mult $_{2}$ and mult $_{4}$ are CFA's. If there exists some CFA that maps $A$ to $B$, then we say that $A$ and $B$ are compatible. Similarly, if there exist some CFA that maps the path/cycle $A_{1}, A_{2}, . ., A_{m}$ to $B_{1}, B_{2}, . ., B_{m}$ we say that $A_{1}, A_{2}, . ., A_{m}$ and $B_{1}, B_{2}, . ., B_{m}$ are compatible. The set of all vertices which are compatible with $A$ will be referred to as the compatibility class of $A$, denoted $\operatorname{com}(A) . A$ is compatible with itself since the identity


Figure 3.9:
automorphism is a CFA. If $n$ and $k$ are prime, all the compatibility classes, with the exception of the compatibility class that contains only the two central vertices, have $k$ vertices in them.

### 3.7 Catalan Numbers

The Catalan numbers arise naturally in combinatorics in many counting problems; often these problems involve recursively defined objects. The $k^{\text {th }}$ Catalan number, $C_{k}$, is defined as $\frac{1}{k+1}\binom{2 k}{k}$ Dejter noted that the number of vertices in $R_{k}$ is always a Catalan number (specifically $R_{k}$, has $C_{k}$ vertices ). An example of one of the ways that Catalan numbers arise is as the result of counting the number of binary strings of length $2 k$ with $k$ 1's and $k 0$ 's such that no initial segment of the string has more

0's than l's. The number of these string is the Catalan number $C_{k}$. If we add a prefix 1 to such a string, it will have the property that every prefix will have more 1's than 0's. This leads us to the following definition:

Definition 3.7.1 We will refer to a binary string as valid if every prefix of that string (including the prefix which is the entire string) has more 1's than 0's.

Note that we can create $C_{k}$ valid bit strings of length $n$ by taking the $C_{k}$ strings of length $2 k$ that have $k$ 1's and $k 0$ 's and the property that no initial segment of the string has more 0's than 1's and adding the prefix 1 to each of them.

THEOREM 3.7.2 Among all the circular shifts of a bit string with $k+1$ 1's and $k 0$ 's there is exactly one valid bit string.

Proof: Suppose toward a contradiction that there are two distinct valid bitstrings $A$ and $A^{\prime}$ such that $s h^{j}(A)=A^{\prime}$ for some $j$. Since $A=\left\langle a_{1}, a_{2}, \ldots a_{n}\right\rangle$ is valid $\sum_{i=1}^{s}(-1+$ $\left.2 a_{i}\right)>0$ for all $1 \leq s \leq n$. Specifically this means that $\sum_{i=1}^{n+1-j}\left(-1+2 a_{i}\right)>0$. Since $1=\sum_{i=1}^{n+1-j}\left(-1+2 a_{i}\right)+\sum_{i=n+1-j}^{n}\left(-1+2 a_{i}\right)$, we know that $\sum_{i=n+1-j}^{n}\left(-1+2 a_{i}\right) \leq$ 0 . But $s h^{j}(A)=A^{\prime}$ which means that $A^{\prime}=\left\langle a_{1-j}, a_{2-j}, \ldots, a_{n-j}\right\rangle=\left\langle a_{1}^{\prime}, a_{2}^{\prime}, \ldots, a_{n}^{\prime}\right\rangle$. This means that $0 \leq \sum_{i=n+1-j}^{n}\left(-1+2 a_{i}\right)=\sum_{i=1}^{j}\left(-1+2 a_{i}^{\prime}\right)$ and hence the prefix $a_{1}^{\prime}, a_{2}^{\prime}, \ldots, a_{n-j}^{\prime}$, of $A^{\prime}$ does not contain more 0 's than 1 's, a contradiction.

Therefore, among all the circular shifts of of a bit string with $k+1$ 's and $k 0$ 's, there is at most 1 valid bitstring. However, we can create $C_{k}$ valid strings with $k+1$ 1's and $k 0$ 's (as noted above) and there are $C_{k}$ classes of circular shifts. Hence there is exactly one valid bit string among the circular shifts of a bit string with $k+11$ 's and $k 0$ 's.

The above theorem tells us that every vertex in $R_{k}$ has a unique representation as a valid string. This is useful in both building and searching through $R_{k}$. All valid strings of length $n$ can be generated recursively [27]. However, for the purpose of out experiments, we choose to generate them non-recursively.

We know that the first two bits of a valid bit string must both be 1 . We will start there and sequentially add bits to the end of the existing string. Let us say we have some partially completed string $A=\left\langle a_{1}=1, a_{2}=1, a_{3}, \ldots, a_{m}\right\rangle$. We then count the number of ones in the partially completed string of length $m$. If this number is $k+1$, we add $n-m 0$ 's to the end of the string and we have a valid string of length $n$. If this is not the case, we certainly have less then $k+11$ 's, and then we will compute the sum $S=\sum_{i=1}^{m}\left(-1+2 a_{i}\right)$. If this sum is greater than 1 , we replace $A$ with the two partially completed strings $\left\langle a_{1}=1, a_{2}=1, a_{3}, \ldots, a_{m}, a_{m+1}=0\right\rangle$ and $\left\langle a_{1}=1, a_{2}=\right.$ $\left.1, a_{3}, \ldots, a_{m}, a_{m+1}=1\right\rangle$. If the sum is 1 , then adding a 0 to the end of the string would make it non-valid, so we replace $A$ with $\left\langle a_{1}=1, a_{2}=1, a_{3}, \ldots, a_{m}, a_{m+1}=1\right\rangle$. If we continue extending the partial strings in this way we will end up with a list of all valid string of length $n$.

Having a unique representation for each vertex of $R_{k}$ allows us to construct the edges of $R_{k}$ easily. Starting with a valid string $A=\left\langle a_{1}, a_{2}, \ldots a_{n}\right\rangle$, for each bit $a_{i}$ with $a_{i}=1$, we flip every bit except $a_{i}$ to get a bit string $B$, and then, if $B$ is not valid, we can simply look through the shifts of $B$ until we find one that is. If we do this for all possible $a_{i}$ 's we will have all the neighbors of $A$. This is quite useful, because we can now construct $R_{k}$ directly without needing to construct $M_{k}$ first.

## Chapter 4

## Reducing the Middle Levels

## Problem

### 4.1 A Reduction of the Middle Levels Problem

In the previous sections we developed the machinery that we will use in this chapter to propose a method different than those used by previous authors for determining whether a particular graph $M_{k}$ is Hamiltonian. We believe that this method has the potential to verify that a particular $M_{k}$ is Hamiltonian in a more computationally efficient manner. Much like previous methods, our method involves a Hamiltonian cycle/path heuristic. The difference lies in how the problem is broken down into steps.

In Observation 3.4.25, we justified that the reduced graph $R_{k}$ has non-trivial auto-
morphisms. We will utilize these automorphisms in order to create a graph which has fewer vertices than $R_{k}$. This is done in a way similar to the way $R_{k}$ is constructed from $M_{k}$ (discussed earlier), with one major difference. In $R_{k}$, every vertex, represents $2 n$ vertices from $M_{k}$. In this new reduced graph, which we will refer to as $Z_{k}$, each vertex represents $k$ vertices in $R_{k}$, but there are two vertices in $R_{k}$, which are not represented in $Z_{k}$. These are the two central vertices whose existence was shown in Theorem 3.5.9.

There are 3 major steps in our method:

1) We construct $Z_{k}$ for the particular $R_{k}$.
2) We use a Hamiltonian path heuristic to find a Hamiltonian cycle in $Z_{k}$. This cycle corresponds to $k$ disjoint cycles in $R_{k}$.
3) We interconnect the cycles as well as the two central vertices into one Hamiltonian path by removing one edge from each cycle and interconnecting them using edges in $R_{k}$. The Hamiltonian cycle in step 3 corresponds to a Hamiltonian cycle in $M_{k}$.

We would like to remind the readers that our explorations are only for the case where both $n$ and $k$ are prime. If this is not the case, we have no guarantee that the number of central vertices is 2 , or that each compatibility class has $k$ vertices in it. For example, $R_{6}$ has 6 central vertices.

We illustrate this method for the case $k=5$ and $n=11$. (Note that these are both primes):

We start off with $R_{5}$ which is pictured below. The two central vertices are placed
over top of each other to de-clutter the image. They are represented as a star and the 5 central vertices are represented as circles. All other vertices are represented with points.


Figure 4.1:

Now we find $Z_{5}$, which is isomorphic to 5 disjoint subgraphs of $R_{5}$.


Figure 4.2:

We now find a Hamiltonian cycle in $Z_{5}$ which corresponds to 5 cycles in $R_{k}$.


Figure 4.3:

Breaking a single edge in each cycle, we interconnect the broken cycles into a long Hamiltonian path in $R_{k}$.


Figure 4.4:

This Hamiltonian path corresponds to a Hamiltonian Cycle in $M_{k}$, by Déjter's Theorem.

### 4.2 Another Possible Reduction

In the previous section we described a method that at least conceptually reduces part of the problem. The issue with that method is that we never specified how to construct $Z_{k}$ or which vertices specifically were in it. Methods of choosing which vertices are in $Z_{k}$ may be time consuming and inefficient. Depending on which method we choose it may actually be more time consuming to construct $Z_{k}$ than
to just simply use a Hamiltonian path heuristic to search through $R_{k}$. Luckily, we do not need to construct $Z_{k}$ in order to utilize the symmetries of $R_{k}$ to reduce the problem. If we have a cycle in $R_{k}$ which includes exactly one vertex from each non-central compatibility class and includes no central vertices, then we can build $k$ compatible cycles and interconnect them just like in the previous method.

Pósa [20] developed a useful Hamiltonian path heuristic which was used in [25], [24] to construct Hamiltonian paths in $R_{k}$. The heuristic starts off at some vertex, $v_{1}$. This vertex will be at the start of our path. Then the heuristic keeps extending the path by randomly choosing an available neighbor of the vertex at the end of the path. It does this until there are no available neighbors with which to extend the path. If $v_{m}$ is the vertex at the end of the path, $v_{1}, v_{2}, \ldots, v_{m}$ and $v_{m}$ has no available neighbors, then the neighbors of $v_{m}$ must all be on the path already. Say $v_{m}$, is adjacent to $v_{i}$, with $1<i<m-1$. We observe that $v_{1}, v_{2}, v_{3}, . ., v_{i}, v_{m}, v_{m-1}, v_{m-2}, . ., v_{i+1}$ is also a path of length $m$. However, it may be possible to extend this new path, unlike its predecessor, since the new path does not have $v_{m}$ at its end. We can keep doing this until we get stuck and then backtrack and try to extend the path again. This heuristic was used successfully to find a Hamiltonian path in $R_{17}$, which is a graph with $129,644,790$ vertices.


Figure 4.5:

We need to modify this heuristic slightly when we search through $R_{k}$ since we want to utilize the symmetries of $R_{k}$. Working in a case were $n$ and $k$ are both prime we first remove the two central vertices from our search space. Then starting with an $l$-vertex we start extending a path from that vertex. Every time we add a vertex to the path, we remove that vertex as well as all vertices that are compatible with it from the search space. This guarantees that the path does not contain two vertices that are compatible with one another. Unlike in Pósa's algorithm, when we cannot extend the path, we are not guaranteed that the neighbors of the vertex at the end of the path are on the path. However, we are guaranteed that each neighbor is on a path which is equivalent under some CFA. So our current path, say $v_{1}, v_{2}, \ldots, v_{m}$, cannot be extended and the vertex at the end of the path, $v_{m}$, is adjacent to a vertex, $w_{i}$, with $1<i<m-1$, which is on some path, $w_{1}, w_{2}, \ldots, w_{m}$, which is equivalent under some CFA to $v_{1}, v_{2}, \ldots, v_{m}$. Notice now that $w_{1}, w_{2}, w_{3}, . ., w_{i}, v_{m}, v_{m-1}, v_{m-2}, . ., v_{i+1}$ is also a path of length $m$ and it might be possible to extend this path, unlike its predecessor. Also note that this path contains exactly one vertex from each compatibility class that we have already removed and still starts with an $l$-vertex. If successful, this
algorithm is faster than Pósa's algorithm, because it needs to construct a much shorter path.

We modify the heuristic further to make sure the path is a cycle. This can be done by:

1) We can check if the point at the end of the path is adjacent to the $l$-vertex at the beginning of the path. If it is, we have a Hamiltonian cycle in $Z_{k}$, as desired. If this is not the case, we backtrack and try again. The heuristic applied in [25],[24] used this idea to make sure their path had the desired endpoints.
2) When the path contains all vertices we can still use Pósa's technique to change which vertex is at the end of the path, as in the diagram on page 56 . We do this repeatedly until the path has the desired end vertex.


Figure 4.6:

Once we have constructed a Hamiltonian cycle in $Z_{k}$, we proceed as in the previous section to interconnect the $k$ cycles we have created as well as the two central vertices. We do this with a different modification of the Pósa heuristic. We start with $11^{k} 0^{k}$ and break the cycle that contains it into a path which starts with $11^{k} 0^{k}=(k) . \mathrm{We}$ will refer to the vertex at the end of the path as $E$. This is our initial path. Now, if
$E$ has a neighbor, $A$, in a cycle, $C$, that we have not yet added to our current path, then we can break $C$ into a path by removing one of the two edges incident with $A$, creating a path $C^{\prime}$. Then we attach $A$ to $E$, there by extending our current path. If $E$ has no such neighbor then all of its neighbors are on the current path and we use Pósa's method to choose a different path with a different vertex at its end. We do this until we have a path containing all vertices with the endpoints $11^{k} 0^{k}$ and $1(10)^{k}$. Using this method allows us to add many vertices at a time to our current path. It also places fewer restrictions on the order of the vertices in the final path because this does not force them to appear in the same order that they did in the cycles we constructed earlier.

### 4.3 Implementation

We coded the algorithm described in section 4.2. The java code is included in Appendix A. This code successfully constructed a Hamiltonian path in $R_{5}$ (the smallest case it can be tried on) and $R_{11}$, verifying that the algorithm works. The code output for $R_{5}$ is included in Appendix B. The next case where this algorithm can be applied is $R_{23}$ and this case is still open. Due to time and memory constraints we did not try the algorithm on this case. Also, we did not recode the algorithm used in [25] for direct comparison on the same machine. Hence, we do not have direct proof that our algorithm is faster. We believe that is it faster based on the fact that it operates on a much smaller graph than $R_{k}$.

### 4.4 Some Possible Improvements the Algorithm

It may be worth noting that in our runs of the algorithm, no backtracking was used at all. In spite of that, the algorithm still managed to find an appropriate Hamiltonian path each time it was implemented. It is likely that using strategic backtracking in this algorithm, both when constructing and interconnecting the cycles, would considerably improve runtimes. This would involve backtracking whenever several successive rearrangements do not yielded an extendible path. Our current method just keeps rearranging the path until it is possible to extend it. If we treat our representative vertices in $R_{k}$ as binary integers, note that when generating the list of vertices in $R_{k}$ as described in section 3.7, the list is built in order from the smallest to the largest integer. This means that if we are searching for a vertex in that list we can do it using a binary search instead of a linear one which is what was done in our code. In the algorithm used in [25], instead of rearranging the list of vertices which represent the path, every time the path could not be extended a list of pending rearrangements was updated. Using this list all the rearrangements were implemented simultaneously. This lead to a significant improvement in runtime. It seems feasible to do something similar to improve our algorithm. It is our hope that implementing all these improvements in our algorithm will allow it to successfully find a Hamiltonian path in $R_{23}$. We also, note that it seems feasible to run a similar algorithm on cases where $n$ and $k$ are not both prime. We will just have more central vertices to find and keep track of. Hence it may be possible that a modification of our algorithm could be used in establishing that $M_{18}$ is Hamiltonian.

## Chapter 5

## Applications to Gray Codes

We will mention one particular application of Hamiltonian cycles, Gray codes. We start by describing an application based on Gray codes and then we introduce them.

Consider a rotating device split into sections with a sensor on it which sends a binary signal to a computer telling it which section the sensor is currently on and hence the current position of the device. You can think of this device as a drum in a photocopy machine, a magnetic storage device or an axle of a vehicle (abs system); all these are real life examples of devices that use this idea. First, we use binary numbers as in part (a) of the above diagram to number the positions of the device:


Figure 5.1:

| position | binary representation | position | binary representation |
| :---: | :---: | :---: | :---: |
| 0 | 0000 | 8 | 1000 |
| 1 | 0001 | 9 | 1001 |
| 2 | 0010 | 10 | 1010 |
| 3 | 0011 | 11 | 1011 |
| 4 | 0100 | 12 | 1100 |
| 5 | 0101 | 13 | 1101 |
| 6 | 0110 | 14 | 1110 |
| 7 | 0111 | 15 | 1111 |

Table 5.1:

If the device stops spinning and the sensor is in-between positions then the signal sent to the computer will be a mix of the two adjacent strings. In the above list of binary strings, we notice that certain adjacent positions differ significantly. For example, the $0^{\text {th }}$ position (0000) and the last position (1111) are adjacent, so that if
the sensor is in-between those two positions, then it is possible that the signal sent to the computer could be about any one of the 16 positions. So it would be useful to have a circular list of binary strings such that any two strings which are adjacent on the list differ in exactly one position. If the device is labeled in this way, then when the sensor is between positions, the sensor sends a signal which is one of the positions which it is in-between, hence giving the computer a much more accurate information about the position of the device. The following is a list of binary strings with this minimal change property, (this one corresponds to part (b) of the diagram):

| position | binary representation | position | binary representation |
| :---: | :---: | :---: | :---: |
| 0 | 0000 | 8 | 1100 |
| 1 | 0001 | 9 | 1101 |
| 2 | 0011 | 10 | 1111 |
| 3 | 0010 | 11 | 1110 |
| 4 | 0110 | 12 | 1010 |
| 5 | 0111 | 13 | 1011 |
| 6 | 0101 | 14 | 1001 |
| 7 | 0100 | 15 | 1000 |

Table 5.2:

These minimal change listings are commonly referred to as "Gray codes" named after Frank Gray who was one of the first people to research them. It is easy to show that Gray codes exist for all possible lengths of binary string. We will include the proof of this well-known fact, because we will need it for the discussion that follows.

THEOREM 5.1.1 [11] A listing of all binary strings of length $n$ exists such that any two binary strings which are adjacent on the list differ in exactly one position (the first and last entries in the list are also considered adjacent).

Proof: We will proceed by induction on $n$. For $n=1$, the list 0,1 has the desired property. Now, assume inductively that listings with the desired adjacency properties exist for all $n<N-1$. This means that there exists such a list, $L_{N-1}$, of all binary strings of length $N-1$. Now we build $L_{N}$ in the following way: add a zero to the beginning of every element in $L_{N-1}$, this creates a new list $L_{N-1}^{\prime}$. Now take the list of the elements of $L_{N-1}$ in reverse order and add a 1 to the beginning of each of those strings; call this list $L_{N-1}^{\prime \prime}$. Combine $L_{N-1}^{\prime \prime}$ and $L_{N-1}^{\prime}$ into one list which we will call $L_{N}$. Since every $N$-bit binary string can be created by adding either a 0 or 1 to the beginning of one of the strings of $L_{N-1}, L_{N}$ contains all possible strings of length $N$. Notice that all strings of $L_{N-1}^{\prime \prime}$ have the adjacency property and so do the strings in $L_{N-1}^{\prime}$. Since the first string in $L_{N-1}^{\prime}$ and the last string in $L_{N-1}^{\prime \prime}$ differ in only the first bit and the first string in $L_{N-1}^{\prime \prime}$ and the last string in $L_{N-1}^{\prime}$ differ in only the first bit, every pair of adjacent strings on the list, differ in exactly one position, as desired.

The above proof shows us how to construct Gray codes of all different sizes recursively. Starting with $0,1 \rightarrow 00,01,11,10 \rightarrow 000,001,011,010,110,111,101,100$ etc... Gray codes constructed in this way are called binary reflected Gray codes. Although these are useful, sometimes in practice it become useful to put additional restrictions on these lists. For example, note that, in a binary reflected Gray code,
the first bit changes only twice which is less than any other bit whereas the last bit changes most often. There are situations in which that may be undesirable and we want a Gray code in which every bit changes the same number of times. Such balanced Gray codes were shown to exist for all even lengths of a binary string by Wagner and West [28]. For a general binary string some heuristics have been proposed. There are many variations on the theme of creating Gray codes with various different restrictions. Currently, the Middle Levels problem seems to be the most notorious unsolved Gray code problem.

## Chapter 6

## Conclusion

We developed two new characteristics of $l$-vertices (Lemma 3.4.10 and Theorem 3.4.4) as well as defined and characterized central vertices (Note 3.4.18). Then we used these characteristics to verify how many of these vertices are in particular $R_{k}$ 's (Theorem 3.5.9 and Observation 3.5.5), focusing on the case where both $n$ and $k$ are prime. We explored the role of these vertices in the automorphisms of $R_{k}$ (Section 3.6) and used central vertices to define a particular type of automorphisms, CFAs. These CFAs, as well as the work of Pósa and Déjter played a crucial role in the development of the algorithm described in section 4.2. We believe that the algorithm is amenable to confirm more cases as Hamiltonian. It is also our hope that our work on the automorphisms of $R_{k}$ will aid the development of a complete solution to the Middle Levels problem.

## Bibliography

[1] R. Cada, T. Kaiser, M. Rosenfeld and Z. Ryjacek, Hamiltonian decompositions of prisms over cubic graphs, Discrete Mathematics 286(2004), 45-56.
[2] Y. Chen, Kneser graphs are Hamiltonian for $n \geq 3 k$, Journal of Combinatorial Theory, Series B 80(2000), 69-79.
[3] F. Chung, P. Diaconis, R. Graham, Universal cycles for combinatorial structures, Discrete Mathematics 110(1992), 43-59.
[4] I.J. Dejter, Hamilton cycles and quotients of bipartite graphs, Graph Theory with Applications to Algorithms and Computer Science (1985), 189-199.
[5] I.J. Dejter and J. Quintana, Long cycles in revolving door graphs, Congressus Numerantium 60(1987), 163-168.
[6] I.J. Dejter, J. Cordova and J. Quintana, Two Hamilton cycles in bipartite reflective Kneser graphs, Discrete Mathematics 72(1988), 63-70.
[7] I.J. Dejter, J. Cedeno and V. Jauregui, A note on Frucht diagrams, Boolean graphs and Hamilton cycles, Discrete Mathematics 114(1994), 131-135.
[8] G.A. Dirac, Some theorems on abstract graphs, Proceedings of the London Mathematical Society 2(1952), 69-81.
[9] D.A. Duffus, H.A. Kierstead and H.S. Snevily, An explicit 1-factorization in the middle of the boolean lattice, Journal of Combinatorial Theory, Series A 65(1994), 334-342.
[10] L. Euler, Solutio d'une question curieuse qui ne paroit soumise á aucune analyse, Mem. Acad. Sci. Berlin 15(1759), 310-337.
[11] F. Gray, Pulse code communications, U.S. Patent Number 2632058 (March 1953).
[12] I. Havel, Semipaths in directed cubes, Graphs and Other Combinatorial Topics (1983), 101-108.
[13] P. Horak, T. Kaiser, M. Rosenfeld and Z. Ryjacek, The prism over the middle-levels graph is Hamiltonian, Order 22 (2005) 73-81.
[14] B. Jackson, Hamiltonian cycles in regular 2-connected graphs, Journal of Combinatorial Theory, Series $B 29$ (1980), 27-46.
[15] H.A. Kierstead and W.T. Trotter, Explicit matchings in the middle levels of the boolean lattice, Order 5 (1988), 163-171.
[16] J. Moon and L. Moser, On Hamiltonian cycles in bipartite graphs, Israel Journal of Mathematics 1(1963), 163-165.
[17] L. Lovász, Problem 11. In Combinatorial Structures and their Applications, (1970).
[18] O. Ore, A note on Hamiltonian circuits, The American Mathematical Monthly 67(1960), 55.
[19] O. Ore, Hamiltonian connected graphs, Journal de Mathématiques Pures et Appliquées 42(1963), 21-27.
[20] L. Pósa, Hamiltonian circuits in random graphs, Discrete Math. 14(4) (1976), 359-364.
[21] F. Ruskey, Combinatorial Generation (2001).
www.1stworks.com/ref/RuskeyCombGen.pdf
[22] C.D. Savage, Long cycles in the middle two levels of the Boolean lattice, Ars Combinatoria 35(A) (1993), 97-108.
[23] C.D. Savage and P. Winkler, Monotone Gray codes and the middle levels problem, Journal of Combinatorial Theory, Series A, 70(1995), 230-248.
[24] C.D. Savage and I. Shields, A Hamilton path heuristic with applications to the middle two levels problem, Congressus Numerantium 140(1999), 161-178.
[25] C.D. Savage, I. Shields and B. J. Shields, An update on the Middle Levels problem, submitled (Aug 2006).
[26] I. Shields, Hamilton cycle heuristics in hard graphs, Ph.D. thesis, North Carolina State University, Raleigh, (2004).
www.lib.ncsu.edu/theses/available/etd-03142004-013420/
[27] R. Stanley, Enumerative combinatorics vol. 1, (1997).
[28] D.G. Wagner and J. West, Construction of uniform Gray codes, Congressus Numerantium. 80(1991), 217-223.

## Appendix A

## Java code

package midlevel;
import java.util.ArrayList

import java.util.Scanner;
import java.io. IoException;import Java.util.Random;
public class Main \{
private ArrayList<Node> node;
/* this is the same $k$ as in the thesis and is assigned avalue in public MidLevel*/

$$
\text { final private int } K \text {; }
$$

// MAX is equal to $\mathrm{k}+1$ and is assigned a value in public MidLevel final private int MAX;
/* this is is the same $n$ as in the thesis ( $n=2 k+1$ )
and is assigned a value in public MidLevel */
final private int $N$;
/* This is a list of all the elements in the
subgroup of $Z_{-} n \backslash\{0\}$ which has order $k * /$
final private int [] multiplylist;
class Node extends Object \{
private ArrayList<Boolean> cell;
ODverride
public boolean equals(0bject object)
Node node $=$ (Node) object;
final int SIZE $=$ this.cell.size()
for (int $i=0 ; i<S I Z E ; ~ i++)$
if (this.cell.get (i) ! = node.cell.get (i)) return false;

[^0]```
/* Default Constructor initiate a new AmrayList */
Node() {
    this.cell = new ArrayList<Boolean>();
F
/* Clone Constructor */
Node(Node clone) {
    this.cell = new ArrayList<Boolean>()
    for (int 1 = 0; i < clone.length(); 1++)
                this.cell.add(clone.get(i));
}
/* Gparam value to be added to our Node String */
void add(boolean value) {
this.cell.add(value);
}
boolean get(int index) {
    return this.cell.get(index);
}
/* Shift Node one placing */
void shift() {
    final int size = this.cell.size() - 1;
        final boolean temp = this.cell.get(size);
        for (1nt 1 = size; i > 0; 1--)
            this.cell.get(i, this.cell.get(i - 1));
        this.cell.set(0, temp);
}
/* Oparam x is the multiplier */
void multiply(int x) throws IOException {
        if ( }x=0=0\mathrm{ )
        throw new IOException();
        final int size = this.cell.size();
        Node clone = new Node(this);
        int location;
        for (int i = 0; i < size; i++) {
            location = i * x % size;
            this.cell.set(location, clone.cell.get(i))
    }
    this.validify()
}
/* Flip the ones to zeros and zeros to ones */
void compliment() (
    final int SIZE = this.cell.size();
    for (int i = O; i < SIZE; 1++)
                If (this.cell,get(i) == true)
                    this.cell.set(i, false);
            else this.cell.set(i, true)
}
/* Cheks if the number of ones exceeds the number of zeros
in every prefix */
boolean valid() {
    int count = 0
    final int SIZE = this.cell.size();
    for (int i = 0; i < SIZE; i++) {
                if (this.cell.get(i))
                    count++;
            else
            if (count <= 0)
                if (count <= 0)
    }
    return true;
}
/* Converts to vertex to a vertex in
the same class which is valid */
void validify() {
        boolean flag = true;
            if (1this.valid())
```

```
                                    this.shift();
            olse
                flag = false;
    while (flag);
    }
    ArrayList<Node> neighbours() throus IDException {
        Node clone;
        ArrayList<Node> list = neu ArrayList<Node> ();
        final int SIZE = this.cell.size();
        for (int i = 0; i < SIZE; i++) {
            if (this.cell.get(i)) {
                clone = neu Node(this);
                clone.cell.set(i, false);
                clone.compliment();
                clone.validify();
                list.add(clone);
            }
    }
    if (list.size() != MAX)
            throw new IOException();
    return list;
}
Arraylist<Node> friends() throus IOException {
    Node clone;
    ArrayList<Node> list = new ArrayList();
    for (int i : multiplyList) {
            clone = new Node(this);
            clone.multiply(i);
            clone,validify();
            clone.validify();
    F
    return list;
}
/* Oreturn value number of ones in Node String */
int numberOfOnes() {
    int count =0;
        for (int i = 0; i < this.cell.size(); i++)
            if (this.cell.get(i))
                count++;
    return count;
F
/* Greturn value number of zeros in Node String */
int numberDiZeros() {
    int count =0;
    for (int i = 0; i < this.cell.size(); i++)
                if (!this.cell,get(i))
                count++;
    return count;
    F
@Dverride
public String toString() {
    String terp = "".
        String temp = ""; 
                1f (cell.get(i) =m true)
                temp += "1";
                    else
                    temp += "0";
        return temp;
    J
    int length() {
    return this.cell.size();
}
}
public Main(int value) throus IOException {
    //this initializes the variables used to what the user indicated
```

```
    K = value;
    MAX = K + 1;
    N = 2* K + 1;
    node = new ArrayList<Node>();
    multiplyList = nev int [K];
    multiplyList[0] = 4;
    for (int i = 1; i< < K; i++)
        multiplyList[i] = 4* multiplyList[i - 1] % N;
    if (multiplyList[multiplyList.length - 1] != 1)
        throw new IOException();
}
private void nodeBuilder(Node temp) throws IOException {
    /*this method is used to extend the partial strings when
    building the list of vertices in R_k*/
    if (temp. numberOfOnes() == this.MAX) {
        // System.out.printin(this,node.get(this.node.size() - 1));
    } else if (temp.number0f0nes() > this.MAX)
        throw new IOException();
    else {
        if (temp.number0fOnes() - 1 == temp.numberOfZeros()) {
            temp.add(true);
            nodeBuilder(temp);
            } else {
                Node clone = new Node(temp);
                temp.add(false);
                nodeBuilder(temp);
                clone.add(true);
                nodeBuilder(clone);
            }
        }
}
private int find(ArrayList<Node> x, Node y) throws IOException {
    /*this method finds and returns the possition of the given node
    on the given list*/
        if (x.size() = 0)
            throw new IOException();
        for (int i = 0; i < x.size(); i++) {
            if (x.get(i).equals(y))
                return i;
        }
        throw nev IOException();
F
private boolean finder(ArrayList<Node> x, Node y) throws IOException {
    /*This method returns true if the given node is on the given list
    and false otherwise*/
        if (x.size() = = 0)
            throw new IOException();
            for (int i = 0; i < x.size(); i++) ;
                if (x.get(i).equals(y))
                    return true;
        }
        return false;
}
private ArrayList<Node> appender(ArrayList<Node> x, ArrayList<Node> y )f
    /*This combines the two lists given into one list*/
        ArrayList<Node> temp = x;
        for( int i = 0; i< y.size(); i++) temp.add(y.get(i));
            return temp;
}
private ArrayList<Node> cycleCutter(ArrayList<Node> \(x\), int \(y\), boolean direction ) 1
/*This method returns the nodes in a cycle starting as the given index and in the give direction*/
ArrayList<Node> temp \(=\) new ArrayList<Node>();
if( !direction ) (
for (int \(i=0 ; i<x . s i z e() ; i++)\)
```

```
                                    temp.add(x.get((y+i)%x.size()));
        }else{
            for(int i=0; i<x,size();i++)
                temp.add(x.get((y-i+x.size())%x.size()));
    }
    return temp;
}
private ArrayList<Node> flip(ArrayList<Node> list, int x, int y)
    /*This implements the modificed Posa flip to make a cycle */
throws IOException {
    if(y<2)return list;
    ArrayList<Node> temp = new ArrayList<Node>();
    final int multiplier = this.multiplyList[x];
    for (int i = 0; i< y + 1; i++
        temp, add(list,get(i));
    for (int i = list.size() - 1; i > y; i--) {
        list.get(i).multiply(multiplier);
        temp.add(list.get(i));
    }
    return temp
}
private ArrayList<Node> flip(ArrayList<Node> list, int y)
throws IOException {
    *This method implements the other modified Posa flip which
    is used to change the end point of the current path then trying to
    interconnect the cycles */
        if(y<2)return list;
        ArrayList<Node> temp = nev ArrayList<Node>()
        for (int i=0;i< y + 1; i++)
            temp.add(list.get(i));
            for (int i = list.size() - 1; i> y; i--) {
                temp.add(list.get(i));
            }
            return temp;
}
private ArrayListeN
    /*This method suplise the first list given with all the common
    entries from the second list removed*/
        ArrayList<Node> result = new ArrayList<Node>();
        boolean inList = false;
        for (int i=0; i < x.size(); i++) {
            for (int j = 0; j< y.size(); j++) {
                    if (x.get(i).equals(y.get(j))) {
                    inList = true;
                    break;
                    }
            }
            if (!inList)
                result.add(x.get(i));
                inList = false;
            }
            return result;
j
private ArrayList<Node> purge(ArrayList<Node> x, Node y) throws
IOException {
    final int size = x.size();
    for (int i = 0; 1 < size; i++)
                if (x.get(i).equals(y)) {
                    x.remove(i);
                    break;
            }
            return x;
}
private ArrayList<Node> intersection(ArrayList<Node> x, ArrayList<Node> y)
```

```
throws IDException {
    ArrayList<Node> result = new ArrayList<Node>()
    boolean inList = false;
    for (int i = 0; i < x.size(); i++) {
        for (int j = 0; j< y.size(); j++) {
            if (x.get(i).equals(y.get(j))) {
                inList = true;
                break;
                }
        }
        if (inList)
            result.add(x.get(i));
        inList = false;
    }
    return result;
}
private ArrayList<Node> multiList(ArrayList<Node> x, int y)throus
IOException{
    ArrayList<Node> temp= new ArrayList<Node> ();
    Node node= nep Node();
    for(int i=0; i< x.size(); i++){
        node = new Node( x.get(i) );
        node.multiply(y);
        temp.add(node);
        }
        return temp;
}
private Node zeroSpaces(Node node) {
    /* this is used in nodeBuilder to append the appropriate number
    of 0's to a string then it already had k+1 1's*/
        int adjust = this.N - node.length();
        if (adjust == 0)
            return node
        for (int i = 0; i < adjust; i++)
            node add(false);
        return node;
}
private static void print(ArrayList<Node> x) {
    // prints a list of vertices
        System.out.println("Length = " + x.size());
        for (int i=0; i< x.size(); i++) \ell
        System.out.println( x.get(i));
        }
}
/* Entry point for our program. */
public static void main(String[] ignored) throws IOException {
    Scanner in = new Scanner(System.in);
    System.out.print("Enter value of K: ");
    Main simulation = new Main(in.nextInt());
    simulation.run();
}
private void run() throws IOException {
    // Make the list of vertices in R_k
    Node first = new Node();
    first.add(true);
    first.add(true)
    this.nodeBuilder(first);
    // Make the two central vertices, then remove them from the list of
    // vertices R_k
    Node cent1 = new Node();
    for (int i = 0; i < N; i++) {
        cent1.add(false);
    }
    System.out.println("R_k made successfully");
    System.out.println("running Posa path constructor....");
    cent1.cell.set(0, true);
```

```
for (int i : multiplyList)
    cent1.cell.set(i, true)
Node cent2 = new Node(cent1);
cent2.cell.set(0, false);
cent2.compliment();
cent1.validify();
cent2.validify();
ArrayList<Node> cent = new ArrayList<Node>()
cent.add(cent1);
cent.add(cent2);
/* compare is a copy of the list of all the vertices in
R_k 'node' which we will use to check that all the vertices
are on our final Hamiltonian path*/
ArrayList<Node> compare = new ArrayList<Node> (node);
node = purge(node, cent);
// Now we the list that will turn into the cycle in Z_n
ArrayList<Node> cycle = new ArrayList<Node>O;
cycle.add(node.get(node.size() - 2));
cycle.add(node.get(node.size() - 1));
node = purge(node, cycle.get(0).friends());
node = purge(node, cycle.get(1),friends());
// Initialize variables used in main loop
ArrayList<Node> inter \(=\) new ArrayList<Node>();
ArrayList<Node> flipList = new ArrayList<Node>();
int location \(=-1\);
int multiIndex \(=-1\);
Node choosen \(=\) new Node()
Node flipNode \(=\) new Node();
boolean flip = false;
boolean allGood = true
boolean done \(=\) false;
boolean pathNowCycle = false;
Random rand = neu Random();
while (!done) \{
while (!flip) f
inter \(=\) intersection(node, cycle.get(cycle.size() - 1)
.neighbours())
if (inter.size() \(!=0\) )
choosen \(=\) inter.get(rand.nextInt(inter.size()));
cycle.add(choosen).
node \(=\) purge(node, choosen.friends());
\} else \{
if (node.size \()=0\) ) \{
done \(=\) true;
system.out.println("done")
break;
\} else
flip \(=\) true;
\}
\}
\}
if (!done) \(\{\) flipList \(=\) purge(cycle.get (cycle.size() - 1).neighbours(), cycle.get (cycle.size() - 2));
flipNode \(=\) flipList.get(rand.nextInt(flipList.size())) flipList \(=\) intersection(cycle, flipNode.friends()); if (fliplist.size \()=0\) ) \(\}\) //Do nothing
else\{
multiIndex \(=\) find(flipNode.friends(), flipList.get(0));
flipNode \(=\) flipList.get (0)
location = find(cycle, flipNode);
cycle \(=\) flip(cycle, multiIndex, location);
flip = false;
\(\left\{^{\}}\right.\)
\} else \{
/* verify that cycle was built correcrly*/ for (int \(1=0 ; i<\operatorname{cycle} . \operatorname{size}()-1\); \(i++\) )
If (!finder (cycle.get(i), neighbours(), cycle,get (i + 1))) \{ System.out.printin("error in cycle construction"); allGood = false;
```

${ }_{\text {if }}$
if (allgood) $\{$
System.out.println("Path contains all vertices:");
System.out.println("Attempting to make path into cycle...");
do \{
if (!finder(cycle.get(cycle.size) - 1).neighbours(),
cycle.get (0) ) ) \{
//System. out.println("path... to ... cycle");
lipList $=$ purge(cycle.get(cycle, size() - 1)
neighbours(), cycle.get(cycle.size() - 2));
flipNode $=$ flipList.get (rand.nextInt(flipList
.size()));
fliplist $=1$ intersection(cycle, flipNode.friends());
$f$ (flithersecion (cyci
System.out.println("flipList error");
\} else \{
multiIndex $=$ find(flipNode.friends(), flipList
get (0));
flipNode $=$ flipList.get $(0)$;
location $=$ find (cycle, flipNode);
cycle $=$ flip(cycle, multiIndex, location);
${ }^{3}$
pathNowCycle = true;
allGood $=$ true ;
for (int $i=0 ; i<\operatorname{cycle} . s i z e()-1 ; 1++$ ) $\{$
if (Ifinder(cycle.get(i).neighbours(), cycle
get (i + 1)) ) $\{$
System.out
.println("error in cycle construction");
allGood = false;
3
\}
\}
(pathNowCycle);
if (alloood) \{
System.out.println("cycle complete");
\}
\}
\}
\}
System.out.println("Interconnecting borken cycles....");
ArrayList<Node> finalPath $=$ nev ArrayList<Node>();//Hamiltonian path
$/ *$ This is a array that contains all the cycles that are equivalent
Under CFA's as well the last entry on the array is a list of the
two central vertices */
ArrayList<Node>[] CYCLES $=$ new ArrayList [K +1$]$;
CYCLES [0]=cycle;
boolean diraction
int numberOfCyclesIncluded;
for (int $i=0 ; 1<K ; i++)\{$
CYCLES[i+1] $=$ multiList(cycle,this.multiplyList[i]);
\}
CYCLES [K]=cent;
ArrayList<Node> NeighborList $=$ new ArrayList<Node>O);
direction = rand.nextBoolean();
finalPath = appender(finalPath, cycleCutter (cYcLes[0], 1,direction));
numberOfCyclesIncluded=1;
doneffalse;
Node neighbor =new Node();
int index;
dof
NeighborList = purge(finalPath.get(finalpath.size()-1).neighbours(), finalPath);
if (numberOfCyclesIncluded $=K+1$ ) done=true;
else if (NeighborList.size ()$=0$ ) $f$
/*if this is the case all neighbors of the last vertex on
our current path lie on the path so we Posa flip the path*/
NeighborList =finalPath.get(finalPath.size()-1).neighbours();
neighbor $=$ NeighborList. get (rand.nextInt(NeighborList.size());
finalPath = flip(finalPath, find(finalPath, neighbor));

## telsed

$/ *$ If this is the case the last vertex on our current path has a neighbor on a path we haven't added yet so we can extend the current path by*/
//find on cycle list;
neighbor $=$ NeighborList.get(rand.nextInt(NeighborList.size()));
int $j=0$;
dot
if ( finder (CYCLES [j], neighbor) ) \{ Index = find(CYCLES[j], neighbor);
direction $=$ rand.nextBoolean();
finalPath = appender ( finalPath,
cycleCutter(CYCLES[j],index, direction) );
umberorcyclesIncluded++

## break

$\stackrel{r}{j++}$
fuhile(true);
3
\}while(!done);
System.out.println("Path now contains all verties");
System.out.println("flipping path to get correct endpoints..."); boolean correctEndpoints $=$ false;
Node end =new Node();
end. add(true);
for(int $i=0 ; i<K ; i++)\{$
end. add(true);
end.add(false);
\}
dof
if (finalPath.get(finalPath.size()-1).equals(end)) correctEndpoints =true; else\{

NeighborList $=$ finalPath.get (finalPath.size()-1) neighbours ();
neighbor $=$ NeighborList.get(rand.nextInt(NeighborList.size())); finalPath $=$ flip(finalPath, find(finalPath, neighbor));
\}
\}while(!correctEndpoints);
System.out.println("Done making Hamiltonian Cycle in R_k");
/*Check that the final path is infact Hamiltonian by

1) checking that all the vertices are in it
2) all adjacent enties are neighbors*/
if ( purge(compare,finalPath).size() $==0$ \& compare.size() $=$ finalPath.size() )
$\{$
System.out.println("Final path contains all vertices") allGood = true;
for (int $i=0 ; i<f i n a l P a t h, g i z e()-1 ; i++)$
f
if (!finder(finalPath.get (i).neighbours(),
finalPath.get(i + 1)))
\{
System.out.println("error in cycle construction"); allGood = false;
J
f
if( allGood) print(finalPath);
r

## Appendix B

## Code output

$R_{5}:$

| Enter value of K: | 11101110000 |
| :--- | ---: |
| 5 | 11110011000 |
| R_k made successfully | 11100111000 |
| running Posa path constructor.... | 11110001100 |
| done | 11100110010 |
| Path contains all vertices: | 11010101100 |
| Attempting to make path into cycle... | 11011001010 |
| cycle complete | 11011010100 |
| Interconnecting borken cycles.... | 11010110010 |
| Path not contains all verties | 11101001010 |
| flipping path to get correct endpoints... | 11110101000 |
| Done making Hamiltonian Cycle in R_k | 11010111000 |
| Final path contains all vertices | 11100101100 |
| Length $\mathbf{4 2}$ | 11101001100 |
| 11111100000 | 11011100010 |
| 11111010000 | 11101101000 |
| 11011110000 | 1111000100 |
| 11110010010 | 11110110000 |
| 11011010010 | 11110010100 |
| 11010110100 | 11101011000 |
| 11101010010 | 11100110100 |
| 11110001010 | 11011011000 |
| 11101010100 | 11110100100 |
| 11101100010 | 1111000010 |
| 1110100100 | 11110100010 |
| 11011001100 | 11011101000 |
| 1101100100 | 11100101010 |
| 11111001000 | 11010101010 |

$R_{11}$

## Appendix C

## Diagrams of $M_{k}$ 's and $R_{k}$ 's

C. $1 \quad R_{k}$

All $l$-vertices are pictured with loops. Central vertices are pictured as empty circles.
The remaining vertices are pictured as filled circles.
$R_{1}$ :

$R_{2}$ : Note that this is the only case were there are vertices that are both central and l-vertices.

$R_{3}$ :

$R_{4}$ :

C. $2 M_{k}$

Because $M_{k}$ is vertex transitive we choose not to label $l$-vertices or central vertices any differently than other vertices.
$M_{1}$ :

$M_{2}$ :

$M_{3}$ :

$M_{4}$ :



[^0]:    return true

