# Combinatorial Designs with Certain Inner Structures and Number Theoretic Approaches to Their Existence 

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## Chapter 1

## Introduction

Combinatorial design theory is a branch of combinatorics studying the systems of finite or discrete objects whose arrangements satisfy specified criteria, such as the properties of balance and symmetry. The study of design theory mainly involves the problems of finding a finite set system with restrictions on the membership and intersections, such as block designs and combinatorial codes. On the other hand, it could also involve the spatial arrangements of entries in arrays, such as magic squares and Latin squares.

Combinatorial inventions on magic squares can be traced back to high antiquity in early China. The Luoshu (Luo River Writing) square, which is legendarily believed to have been created between the 3 rd and the 2nd millennium BC , is the earliest record of a $3 \times 3$ magic square with the numbers 1 to 9 . In modern mathematics, the study of design theory has its roots in the work of L. P. Euler who posed the "36 officer problem" in 1782. This problem is equivalent to finding "mutually orthogonal Latin squares (MOLSs)" of order 6. Euler also conjectured that MOLSs of order $n$ do not exist for any $n \equiv 2(\bmod 4)$, which was known as Euler's conjecture until it was shown to be false by Bose, Shrikhande, and Parker [11] in 1960. Later in the 19th century, combinatorial designs were studied as geometric configurations by T. P. Kirkman, J. Steiner, and A. Cayley.

In the 1930s, the development of statistical experimental designs greatly promoted research on combinatorial designs. A milestone of design theory was established by Fisher and Yates, who made use of "balanced incomplete block designs" (BIBDs) and "lattice squares" for agricultural experiments [45, 127, 128. Around the same time, Bose [10] published a long paper on the constructions of BIBDs, in which a most significant construction technique for designs, called "the method of differences", was proposed. Along this direction, the existence and construction of "difference families" has become a rapidly expanding subject with fundamental importance. In terms of contemporary design theory, a BIBD is a 2 -design, and a difference family generates a 2-design with an automorphism group acting sharply transitively on its points. The main problems for design theory in recent stages can be summarized as "existence", "construction", and
"characterization".
Constructions of designs fall into two categories, direct constructions and recursive constructions. For 2-designs, difference families are commonly used for direct constructions, and related studies have been investigated by intensive use of finite fields, finite groups, and algebraic number theory. However, recursive constructions are usually purely combinatorial. Infinite families of new designs can be obtained from known designs via recursive constructions, like an interlocking puzzle game. The most outstanding recursive construction of 2-designs is due to Wilson [121, who proposed the "pairwise balanced design (PBD) construction" and then revolutionarily established the existence of 2designs. Thereafter, using similar ideas, various recursive constructions have been studied, and excellent progress has been made on the existence of designs.

In contrast, only a few constructions for $t$-designs with $t \geq 3$ are known. In the mid 20th century, before the completion of the classification of finite simple groups, group theorists found a close relationship between high transitivity of finite groups and designs. They tried to develop research on finite groups via the study of $t$-designs with large $t[124$. However, non-trivial $t$-designs have been proved to exist for any positive integer $t$ by Teirlinck [113], which leads to a gap with highly transitive finite groups. Recently, Keevash 63] settled the existence of $t$-designs for all but finitely many admissible parameters by a new probabilistic method referred to as "randomized algebraic construction". Nevertheless, to find an explicit construction for a $t$-design with given parameters is still difficult in general. In particular, for the cases when $t \geq 3$, the existence of a $t$-design which is invariant under a prescribed permutation group without high transitivity is still often unknown.

Characterizations of designs study the inner structures of a design, which roughly fall into two aspects, algebraic characterizations and geometric characterizations. For the algebraic aspects, the "automorphism groups" of a design is the main theme. In other words, it is desired to consider designs which are invariant under a permutation group. In particular, cyclic designs are desired for applications to communication systems. For the geometric aspects, a design can be viewed as a "geometry" consisting of "points" and "lines", and the intersections between "lines" are the most important. In particular, if a design can be partitioned into subsystems, each of which forms a "parallel class" of lines, then the design is said to be resolvable. For applications to experiments and group testings, resolvability plays an important role.

In this dissertation, we will concentrate on two kinds of designs, both of which have specific inner structures. Firstly, we consider 3-designs which are invariant under the generalized affine groups, which are a special type of cyclic 3-designs. Secondly, we will study the difference families with respect to "gridblock designs", which are known to play an essential role for DNA library screening and other group testing models for experiments. For both types, approaches to the direct constructions and existence will be given by employing group theoretic and number theoretic tools. Recursive constructions and computational results will also be used to establish their existence.

This chapter is devoted to providing a brief introduction and preliminaries in
design theory. Some basic concepts and properties on $t$-designs will be presented in Sections 1.1 and 1.2 . Moreover, we will focus mainly on 3-designs with prescribed automorphism groups, difference families, and grid-block designs in Sections 1.3, 1.4 and 1.5 respectively.

Before proceeding further, we introduce some notation that will be used throughout this dissertation.

Let $V$ be a finite set and let $X$ be a subset of $V$ with $|X|=k$, where $|X|$ is the cardinality of $X$. Then we also say $X$ is a $k$-subset of $V$ and denote all the $k$-subsets of $V$ by $\binom{V}{k}$. Moreover, let $2^{V}$ denote the set consisting of all the subsets of $V$. We also use $\# X$ to denote the cardinality of $X$, especially when $X$ has a complicated set-builder notation with a vertical bar in it.

Moreover, suppose $\Omega$ is a permutation group acting on $V$. Then $\Omega$ acts naturally on $\binom{V}{k}$ for any positive integer $k$. For any $B \in\binom{V}{k}$, let $\mathcal{O}_{\Omega}(B)$ denote the orbit of $B$ under the action of $\Omega$, that is $\mathcal{O}_{\Omega}(B)=\left\{B^{\omega} \mid \omega \in \Omega\right\}$. Moreover, for any $\mathcal{B} \subseteq 2^{V}$, let $\mathcal{O}_{\Omega}(\mathcal{B})=\bigcup_{B \in \mathcal{B}} \mathcal{O}_{\Omega}(B)$. If $\mathcal{O}_{\Omega}(\mathcal{B})=\mathcal{B}$, then $\mathcal{B}$ is said to be invariant under the action of $\Omega$, or $\Omega$ leaves $\mathcal{B}$ invariant.

Let $G$ be an additive group and let $X, Y$ be subsets of $G$. We denote $X+Y=\{x+y \mid x \in X, y \in Y\}$ and $X+a=\{x+a \mid x \in X\}$ for any $a \in G$.

For positive integers $a, b$ with $a<b$, let $[a, b]$ denote the set $\{a, a+1, \ldots, b\}$. In particular, let $[n]$ denote the set $\{1,2, \ldots, n\}$ for a positive integer $n$. For a real number $x$, let $\lfloor x\rfloor$ denote the the largest integer less than or equal to $x$.

Let $\mathbb{Z}_{n}$ denote the ring of integers modulo $n$, that is $\mathbb{Z} / n \mathbb{Z}$. Let $\mathbb{Z}_{n}^{\times}$and $\mathbb{Z}_{n}^{*}$ denote the multiplicative group and the set of all nonzero elements, respectively, of $\mathbb{Z}_{n}$. We also use $\mathbb{Z}_{n}$ for the cyclic group of order $n$. Let $\mathbb{F}_{q}$ denote the finite field of order $q$ and let $\mathbb{F}_{q}^{*}$ denote the multiplicative group of $\mathbb{F}_{q}$, that is, the set of all nonzero elements of $\mathbb{F}_{q}$. More notation for further discussion over $\mathbb{F}_{q}$ will be introduced in Section 1.4 .

### 1.1 Combinatorial $t$-designs

In this section, we give a brief overview of some basic concepts and important results of combinatorial $t$-designs.

Definition 1.1.1 ( $t$-design). Let $V$ be a finite set of $v$ points, and let $\mathcal{B}$ be a collection of $k$-subsets (blocks) of $V$. The pair $(V, \mathcal{B})$ is called a $t-(v, k, \lambda)$ design if every $t$-subset appears in exactly $\lambda$ blocks of $\mathcal{B}$.

The parameters $k$ and $\lambda$ are called block size and index, respectively. $(V, \mathcal{B})$ is said to be simple if there are no repeated blocks in $\mathcal{B}$. A 2-design is well known as a balanced incomplete block design (BIBD), which is commonly used for experimental designs. In the case of $\lambda=1, t$-designs are also called Steiner systems. In particular, 2- $(v, 3,1)$ designs and $3-(v, 4,1)$ designs are known as Steiner triple systems and Steiner quadruple systems, and denoted by STS $(v)$ and $\operatorname{SQS}(v)$, respectively.

The earliest study involving $t$-designs can be traced to Plücker 95, in 1835, who mentioned a $2-(9,3,1)$ design $(\operatorname{STS}(9))$ in his work on algebraic curves.

Later, in 1839, Plücker [96] described a 3-(28, 4, 1) design (SQS(28)) and then he asked what kind of parameters $t$ and $v$ are realizable for a $t-(v, t+1,1)$ design. In 1844, Woolhouse [125 presented a more general problem: does there exist $(V, \mathcal{B})$ such that any $t$-subset of $V$ is contained in at most one block? Such a design is now known as a $t$-packing design (see Definition 1.2.9).

By counting the number of blocks containing a fixed $i$-subset $I$ for $0 \leq i \leq t$, we have

$$
\#\{B \in \mathcal{B} \mid I \subset B\}=\lambda \frac{\binom{v-i}{t-i}}{\binom{k-i}{t-i}},
$$

which implies the following divisibility conditions:
Proposition 1.1.2 (divisibility conditions). If there exists a $t-(v, k, \lambda)$ design, then

$$
\begin{equation*}
\lambda\binom{v-i}{t-i} \equiv 0 \quad\left(\bmod \binom{k-i}{t-i}\right) \quad \text { for any } 0 \leq i \leq t \tag{1.1}
\end{equation*}
$$

Proposition 1.1.3 (divisibility conditions for 2-designs). If there exists a 2 $(v, k, \lambda)$ design, then

$$
\begin{equation*}
\lambda v(v-1) \equiv 0 \quad(\bmod k(k-1)) \quad \text { and } \quad \lambda(v-1) \equiv 0 \quad(\bmod k-1) \tag{1.2}
\end{equation*}
$$

For Steiner triple systems, Kirkman [64] proved that the divisibility conditions (1.2) are also sufficient. Six years later, without knowing Kirkman's work, Steiner [109] noticed the divisibility condition, that is $v \equiv 1,3(\bmod 6)$, and presented the problem of the existence of Steiner triple systems. Actually, for 2-designs with larger $k$, the results are also plentiful. The existence problem of $2-(v, 4,1)$ designs and $2-(v, 5,1)$ designs were solved by Hanani 53, 54] in 1961 and 1972, respectively. For $6 \leq k \leq 9$, the existence of 2-designs are nearly completely settled except for some small parameters (see $34 \S 3.1$ for details). It is notable that, in order to show the existence of 2-designs, Wilson $121,122,123$. developed an outstanding construction, called PBD-construction, and finally proved the following theorem:

Theorem 1.1.4 (Wilson [123]). For given $k$ and $\lambda$ and for a sufficiently large $v \geq v_{0}(k, \lambda)$, the divisibility conditions (1.2) are sufficient for the existence of a $2-(v, k, \lambda)$ design.

It is nature to consider a $2-(v, k, \lambda)$ design as a decomposition of $K_{v}^{(\lambda)}$ into $k$-cliques, where $K_{v}^{(\lambda)}$ denotes the complete multi-graph with $v$ vertices such that there are $\lambda$ edges between every pair of vertices. In particular, when $\lambda=1$, $K_{v}^{(1)}$ is the complete graph of order $v$ in the usual sense. We can also replace " $k$ clique" to any other subgraph of $K_{v}$ for "graph decompositions" in general. For "graph decompositions", there are plenty of literature, and the formal definitions will be given in Section 1.5 . For more details, the interested readers may refer to 34 §VI.24.

In contrast, the constructions for $t$-designs with $t \geq 3$ are quite rare. First, we give a brief historical review on the study of Steiner quadruple systems,
which are the minimal nontrivial classes for $t \geq 3$. In 1847, Kirkman 64] proved that for any positive integer $n$, there exists an $\operatorname{SQS}\left(2^{n}\right)$. Nearly 70 years later, Fitting [46 constructed an $\mathrm{SQS}(26)$ and an $\mathrm{SQS}(34)$ by proposing a graph theoretic construction. Eventually, the existence of Steiner quadruple systems was settled by Hanani [52] in 1960 via a series of complicated recursive constructions. Thereafter, Lenz [70] and Hartman [55] simplified Hanani's proof. Recently, Zhang and Ge [132] proposed a new proof which is more elegant and more concise.

Theorem 1.1.5 (Hanani [52]). There is an $S Q S(v)$ if, and only if, $v \equiv 2,4$ $(\bmod 6)$.

The first existence result of $t$-designs dealing with general $t$ is due to Teirlinck [113] who showed the following theorem:

Theorem 1.1.6 (Teirlinck [113]). There exists a non-trivial t-design for any positive integer $t$.

It should be mentioned that Keevash 63 posted work in 2014 which claims that the divisibility conditions 1.1 for a $t-(v, k, \lambda)$ design are sufficient for all but finitely many admissible parameters. Keevash's proof relies deeply on probabilistic combinatorics, and his new method is referred to as "randomized algebraic construction". Nevertheless, the constructions for a $t$ - $(v, k, \lambda)$ design are still of interest, from both theoretical and application reasons.

Finally, we introduce the concept of resolvability for $t$-designs. A $t$-design $(V, \mathcal{B})$ is said to be resolvable if $\mathcal{B}$ can be partitioned into parallel classes (also known as resolution classes), where each parallel class is a partition of $V$.

### 1.2 Automorphism groups of $t$-designs and applications to optical communications

In this section, we focus on $t$-designs with specific algebraic inner structures, that is, $t$-designs admitting specific automorphism groups. Also, some applications of cyclic designs to optical communications will be interpreted.

Definition 1.2 .1 (automorphism groups of $t$-designs). Let $G$ be a permutation group acting on $V$ and let $(V, \mathcal{B})$ be a $t-(v, k, \lambda)$ design. If $G$ leaves $\mathcal{B}$ invariant, then $G$ is called an automorphism group of $(V, \mathcal{B})$. In this case, $\mathcal{B}$ can be partitioned into orbits under the action of $G$. We can choose any block in an orbit as a base block to represent the whole orbit. For any $B \in \mathcal{B}$, if $\left|\mathcal{O}_{G}(B)\right|=$ $|G|$, then $\mathcal{O}_{G}(B)$ is said to be full, otherwise short.

Furthermore, if $|V|=v$ and $G$ is the cyclic group of order $v$, then the orbits are called cyclic orbits and $(V, \mathcal{B})$ is said to be cyclic. If $\mathcal{B}$ gives no short orbit under the action of $G$, then $(V, \mathcal{B})$ is said to be strictly cyclic. In these cases, the point set $V$ can be identified with $\mathbb{Z}_{v}$. Strictly cyclic 2 -designs are equivalent to cyclic difference families, which will be discussed in detail in Section 1.4.

Example 1.2.2. Let $V=\mathbb{Z}_{10}$ and let $\mathcal{B}$ be the collection of the following blocks:

| $\{0,1,5,9\}$, | $\{0,2,5,8\}$, | $\{0,1,3,4\}$, |
| :--- | :--- | :--- |
| $\{1,2,6,0\}$, | $\{1,3,6,9\}$, | $\{1,2,4,5\}$, |
| $\{2,3,7,1\}$, | $\{2,4,7,0\}$, | $\{2,3,5,6\}$, |
| $\{3,4,8,2\}$, | $\{3,5,8,1\}$, | $\{3,4,6,7\}$, |
| $\{4,5,9,3\}$, | $\{4,6,9,2\}$, | $\{4,5,7,8\}$, |
| $\{5,6,0,4\}$, | $\{5,7,0,3\}$, | $\{5,6,8,9\}$, |
| $\{6,7,1,5\}$, | $\{6,8,1,4\}$, | $\{6,7,9,0\}$, |
| $\{7,8,2,6\}$, | $\{7,9,2,5\}$, | $\{7,8,0,1\}$, |
| $\{8,9,3,7\}$, | $\{8,0,3,6\}$, | $\{8,9,1,2\}$, |
| $\{9,0,4,8\}$, | $\{9,1,4,7\}$, | $\{9,0,2,3\}$. |

Then $(V, \mathcal{B})$ is a $3-(10,4,1)$ design $(\mathrm{SQS}(10))$, which is strictly cyclic. Each cyclic orbit consists of the ten blocks listed in each column.

Then, we introduce the concept of the affine-invariant property for a cyclic $t$-design.

Definition 1.2.3 (multiplier). Let $\left(\mathbb{Z}_{v}, \mathcal{B}\right)$ be a cyclic $t$-design and let $\alpha$ be a unit in $\mathbb{Z}_{v}$, that is $\alpha \in \mathbb{Z}_{v}^{\times}$. For any $B \in \mathcal{B}$, if $\alpha B \in \mathcal{B}$, then $\alpha$ is called a multiplier of $\left(\mathbb{Z}_{v}, \mathcal{B}\right)$.

Definition 1.2.4 (affine-invariant $t$-design). A cyclic $t$-design $\left(\mathbb{Z}_{v}, \mathcal{B}\right)$ is said to be affine-invariant, if every $\alpha \in \mathbb{Z}_{v}^{\times}$is a multiplier.

In other words, an affine-invariant $t$-design $\left(\mathbb{Z}_{v}, \mathcal{B}\right)$ admits the group $A$ as its automorphism group, where $A$ is the general affine group of degree one over $\mathbb{Z}_{v}$ defined by

$$
A=\left\{(i, \alpha) \mid i \in \mathbb{Z}_{v}, \alpha \in \mathbb{Z}_{v}^{\times}\right\} \cong \mathbb{Z}_{v} \rtimes \mathbb{Z}_{v}^{\times}
$$

For a subset $S \subseteq \mathbb{Z}_{v}$, the affine orbit of $S$ is the orbit of $S$ under the action of the general affine group $A$, denoted by $\mathcal{O}_{A}(S)$.

Example 1.2.5. The unique (up to isomorphism) SQS(10) in Example 1.2 .2 is affine-invariant. Take

$$
B_{1}=\{0,1,5,9\}, \quad B_{2}=\{0,2,5,8\}, \quad \text { and } B_{3}=\{0,1,3,4\}
$$

as base blocks of the cyclic orbits. We have $B_{1} \times 3+5=\{0,3,5,7\}+5=$ $\{5,8,0,2\}=B_{2}$ over $\mathbb{Z}_{10}$. Hence, the cyclic orbits of $B_{1}$ and $B_{2}$ are contained in the same affine orbit. In fact, there are only two affine orbits, namely, $\mathcal{O}_{A}\left(B_{1}\right)$ $\left(=\mathcal{O}_{A}\left(B_{2}\right)\right)$ and $\mathcal{O}_{A}\left(B_{3}\right)$.

Now we begin to consider a family of combinatorial codes. An optical orthogonal code (briefly, OOC) is a binary code with good auto- and cross-correlation properties. The study of OOCs is motivated by applications to code-division
multiple-access (CDMA) communication by fiber-optic channels. Low auto- and cross-correlations can efficiently reduce conflicts with undesired signals in communication (see [32]). We first give the definition of optical orthogonal codes from the point of view of binary codes.

Definition 1.2.6 (optical orthogonal code). A binary code $\mathcal{C} \subseteq\{0,1\}^{n}$ is called an optical orthogonal code with parameter $\left(n, k, \lambda_{a}, \lambda_{c}\right)$ (briefly, an ( $n, k, \lambda_{a}, \lambda_{c}$ )OOC), if the following properties hold:
(i) For any codeword $\boldsymbol{x}=\left(x_{0}, x_{1}, \cdots, x_{n-1}\right) \in \mathcal{C}$, the Hamming weight $\mathrm{w}(\boldsymbol{x})=k ;$
(ii) For any codeword $\boldsymbol{x}=\left(x_{0}, x_{1}, \cdots, x_{n-1}\right) \in \mathcal{C}$ and any relative delay offset $\tau \not \equiv 0(\bmod n)$, the Hamming auto-correlation of $\boldsymbol{x}$ satisfies

$$
H_{\boldsymbol{x}}(\tau)=\sum_{t=0}^{n-1} x_{t} x_{t+\tau} \leq \lambda_{a}
$$

(iii) For any pair of distinct codewords $\boldsymbol{x}=\left(x_{0}, x_{1}, \cdots, x_{n-1}\right) \in \mathcal{C}, \boldsymbol{y}=$ $\left(y_{0}, y_{1}, \cdots, y_{n-1}\right) \in \mathcal{C}$ and any relative delay offset $\tau$, the Hamming crosscorrelation of $\boldsymbol{x}$ and $\boldsymbol{y}$ satisfies

$$
H_{\boldsymbol{x}, \boldsymbol{y}}(\tau)=\sum_{t=0}^{n-1} x_{t} y_{t+\tau} \leq \lambda_{c}
$$

where the subscripts of $x_{i}$ and $y_{j}$ are reduced modulo $n$. In particular, if $\lambda_{a}=$ $\lambda_{c}=\lambda$, we simply write an $(n, k, \lambda)$-OOC. The number of codewords in $\mathcal{C}$ is called the size of $\mathcal{C}$.

Alternatively, the definition can be rephrased by considering the support of each codeword $\boldsymbol{x}$, that is the set of indices of all nonzero coordinates of $\boldsymbol{x}$. Then we can identify $\mathcal{C}$ with a collection of $k$-subsets of $\mathbb{Z}_{n}$.

Definition 1.2.7 (an optical orthogonal code as a set system). Let $\mathcal{C}$ be a collection of subsets of $\mathbb{Z}_{n}, \mathcal{C}$ is called an $\left(n, k, \lambda_{a}, \lambda_{c}\right)$ optical orthogonal code if
(i)' For any $X \in \mathcal{C},|X|=k$;
(ii)' For any $X \in \mathcal{C}$ and any nonzero $\tau \in \mathbb{Z}_{n},|X \cap(X+\tau)| \leq \lambda_{a}$;
(iii) ${ }^{\prime}$ For any distinct $X, Y \in \mathcal{C}$ and any $\tau \in \mathbb{Z}_{n},|X \cap(Y+\tau)| \leq \lambda_{c}$.

For given parameters $(n, k, \lambda)$, we denote the largest possible size of an $(n, k, \lambda)$-OOC by $\Phi(n, k, \lambda)$. An $(n, k, \lambda)$-OOC $\mathcal{C}$ is said to be optimal if $|\mathcal{C}|=$ $\Phi(n, k, \lambda)$. In general, it is hard to determine the exact value of $\Phi(n, k, \lambda)$ for certain $(n, k, \lambda)$. However, an OOC can be seen as a constant weight errorcorrecting code. Thus $\Phi(n, k, \lambda)$ can be bounded above by the Johnson bound $J(v, k, \lambda)$ 62].

Proposition 1.2.8 (Johnson bound 62]).

$$
\Phi(n, k, \lambda) \leq\left\lfloor\frac{1}{k}\left\lfloor\frac{n-1}{k-1}\left\lfloor\frac{n-2}{k-2}\left\lfloor\cdots\left\lfloor\frac{n-\lambda}{k-\lambda}\right\rfloor \cdots\right\rfloor\right\rfloor\right\rfloor\right\rfloor=: J(v, k, \lambda)
$$

Optimal OOCs are closely related to $t$-designs and $t$-packings.
Definition 1.2.9. A $t-(v, k, \lambda)$ packing design, or simply a $t$ - $(v, k, \lambda)$ packing, is a pair $(V, \mathcal{B})$, where $V$ is a set of $v$ points, and $\mathcal{B}$ is a collection of $k$-subsets (blocks) of $V$, such that any $t$-subset of $V$ appears in at most $\lambda$ blocks.

Clearly, a $t$-packing generalizes the notion of a $t$-design. Similarly, a $t$ $(v, k, \lambda)$ packing is cyclic if it admits $\mathbb{Z}_{v}$ as its automorphism group, and it is strictly cyclic if all the cyclic orbits are full. A strictly cyclic $t-(v, k, \lambda)$ packing is said to be optimal if the number of base blocks, say $b$, attains the following equality (see [102]):

$$
\begin{equation*}
\left.b \leq\left\lfloor\frac{1}{k}\left\lfloor\frac{v-1}{k-1}\left\lfloor\frac{v-2}{k-2}\left\lfloor\cdots\left\lfloor\frac{v-t+1}{k-t+1} \lambda\right\rfloor \cdots\right\rfloor\right\rfloor\right\rfloor\right\rfloor\right\rfloor \tag{1.3}
\end{equation*}
$$

There is an equivalence between an optimal OOC and an optimal $t$-packing.
Theorem 1.2.10 (Fuji-Hara and Miao [48, Chu [31). Any optimal ( $v, k, \lambda)$ OOC is equivalent to an optimal strictly cyclic $(\lambda+1)-(v, k, 1)$ packing, provided that $1 \leq \lambda<k-1$.

In particular, a strictly cyclic $\operatorname{SQS}(v)$ is equivalent to an optimal $(v, 4,2)$ OOC. In Section 1.3 , we will introduce some important results on strictly cyclic SQSs.

### 1.3 Quadruple systems with a prescribed automorphism group

In this section, we draw our attention to quadruple systems, that is $3-(v, 4, \lambda)$ designs. As stated before, we write an $\operatorname{SQS}(v)$ for a $3-(v, 4,1)$ design. Moreover, we write a $\operatorname{TQS}(v)$ for a two-fold quadruple system, that is a $3-(v, 4,2)$ design. Firstly, we will briefly introduce the results on cyclic SQSs. Let $\sigma$ denote the permutation on $\mathbb{Z}_{v}$ defined by $a^{\sigma}=-a$, and let $\widehat{\mathbb{Z}_{v}}=\mathbb{Z}_{v} \rtimes\langle\sigma\rangle$. A block $B$ is said to be symmetric if $\mathcal{O}_{\mathbb{Z}_{v}}(B)=\mathcal{O}_{\widehat{\mathbb{Z}_{v}}}(B)$. A cyclic $\operatorname{SQS}(V, \mathcal{B})$ is said to be symmetric if every block in $\mathcal{B}$ is symmetric and $\mathcal{B}$ is invariant under the action of $\widehat{\mathbb{Z}_{v}}$.

Now we denote a cyclic SQS by a $C S Q S$, a strictly cyclic SQS by an $s S Q S$, and a symmetric cyclic SQS by an $S$-cyclic $S Q S$. A divisibility condition for the existence of an $\operatorname{sSQS}(v)$ can be easily shown.

Proposition 1.3.1 (Köhler [65]). If an $\operatorname{sSQS}(v)$ exists, then $v \equiv 2,10(\bmod 24)$.

We have mentioned that Fitting [46] constructed an SQS(26) and an SQS(34) in 1915. In fact, those SQSs are S-cyclic. The most famous construction of Scyclic sSQSs is due to Köhler [65], who introduced the notion of "first Köhler graphs".

Theorem 1.3.2 (Köhler [65]). Let $v \equiv 2,10(\bmod 24)$. If the first Köhler graph has a 1-factor, then there exists an sSQS(v).

In order to facilitate the process of finding a 1-factor, Köhler 65 investigated the multiplier automorphisms on those graphs, and introduced the "Köhler orbit graphs" (see also [12, 66, 71]). Along this direction, Siemon [103] investigated "Köhler orbit graphs" for a few more parameters.

Lemma 1.3.3 (Köhler [65, Siemon [103]). There exists an $s S Q S(v)$ if
(i) $v \in\{2,10,26,34,50,58,74,82,106,178,202,226,274,298,346,394$, $466,586,634\}$,
(ii) $v \in\{122,170,194,314,338,386,458,578\}$.

By observing the "embedding" structure of "Köhler orbit graphs", Siemon [103, 104 proposed infinite families of sSQSs.

Theorem 1.3.4 (Siemon 103, 104). Let $m$ be a positive integer. For $p \equiv 5$ $(\bmod 12)$, if the "Köhler orbit graph" with respect to an $s S Q S(2 p)$ has a 1-factor, then an sSQS $\left(2 p^{m}\right)$ exists.

Thereafter, Siemon 105 found that the existence of 1-factors of "Köhler orbit graphs" can be reduced to a number theoretic conjecture called "complete interval conjecture" (see also [2] Problem 146), and verified the conjecture for more parameters.

Theorem 1.3.5 (Siemon 105$)$. There exists an $S$-cyclic $s S Q S(2 p)$ for all prime $p \equiv 53,77(\bmod 120)$ and $p<500000$.

Piotrowski presented a number of important results on S-cyclic sSQSs in his dissertation [93] (see also [104]).

Theorem 1.3.6 (Piotrowski [93]). There exists an $S$-cyclic $s S Q S(2 p)$ for a prime $p$ if
(i) $p \equiv 1(\bmod 4)$ and $p \leq 229$, or
(ii) $p \equiv 1(\bmod 4)$ and $p \not \equiv 1,49(\bmod 60)$ and $p<15000$.

Theorem 1.3.7 (Piotrowski 93 Satz 14.1). There exists an S-cyclic $S Q S(v)$ if and only if $v \equiv 0(\bmod 2), v \not \equiv 0(\bmod 3), v \not \equiv 0(\bmod 8), v \geq 4$, and there exists an $S$-cyclic $S Q S(2 p)$ for any prime divisor $p$ of $v$.

Bitan and Etzion [8] extended Köhler's graph construction 65] and refined Siemon's "complete interval conjecture" 105] for S-cyclic SQS(4p). They verified the number theoretic conjecture by computer programs and showed the following:

Theorem 1.3.8 (Bitan and Etzion [8]). There is an S-cyclic $S Q S(4 p)$ for any prime $p \equiv 5(\bmod 12)$ with $p<1500000$.

For more information about CSQSs and SQSs with other specified automorphism groups, the reader may refer to Lindner and Rosa 78, Grannel and Griggs [51, Hartman and Phelps [55], and Siemon [107]. Notably, Munemasa and Sawa 83 generalized "Köhler orbit graphs" and Piotrowski's theorem on CSQSs to an abelian group $A$ whose Sylow 2-subgroup is cyclic, and established the theory for symmetric $A$-invariant SQSs.

Recursive constructions of 3-designs with a point-regular automorphism group are more complicated than that of 2 -designs (difference families). For recent progress on recursive constructions, the reader may refer to Feng, Chang, and Ji 44], Feng and Chang [43], and Li and Ji [72].

Next, we consider an affine-invariant $\operatorname{sSQS}(v)$ which is simply written as an $\operatorname{AsSQS}(v)$. In general, the number of affine orbits is much less than that of cyclic orbits. For instance, the number of affine orbits of the $\operatorname{AsSQS}(v)$ obtained by our Construction 2.2 .20 is approximately $\frac{1}{6} v$. In contrast, the number of cyclic orbits is approximately $\frac{1}{24} v^{2}$ (see Tables 2.4 and 2.5 in Section 2.2.3. From the viewpoint of OOCs , the storage requirements of codewords are reduced by up to $\frac{1}{4} v$ times. Therefore, we are willing to consider an AsSQS rather than just an sSQS.

Constructions for AsSQSs are less known. Piotrowski 93 proved the following Theorem 1.3.9 (see also [104]). All the previously mentioned results use the aid of some graphs.

Theorem 1.3.9 (Piotrowski [93]). There exists an AsSQS(2p) for prime $p \equiv 1$ $(\bmod 4)$ if $p \not \equiv 1,49(\bmod 60)$ and $p<15000$, or $p \leq 229$.

On the other hand, without the help of graphs, Yoshikawa [129] independently presented an algorithm for constructing an $\operatorname{AsSQS}(2 p)$ and obtained the following theorem:

Theorem 1.3.10 (Yoshikawa [129]). There exists an $A s S Q S(2 p)$ for prime $p \equiv$ $1,5(\bmod 12)$ with $17 \leq p<200$.

Although the parameters are covered by Piotrowski's Theorem 1.3.9, the resulting AsSQSs can be shown to be non-isomorphic. We will characterize Yoshikawa's idea and propose a criterion for the existence of this kind of AsSQSs in Section 2.2.3,

Furthermore, affine-invariant $3-(v, 4, \lambda)$ designs for $\lambda \geq 2$ have also been studied by Köhler [67], who proposed necessary and sufficient conditions for the existence of an affine-invariant $3-(p, 4, \lambda)$ design when $p$ is prime and $\lambda \geq 2$ (see also [13]). However, Köhler's construction for a $\operatorname{TQS}(p)$ also relies on "Köhler orbit graphs", and hence did not able to give an infinite family. In Sections 2.5 and 2.6. we will consider an affine-invariant $\operatorname{TQS}(p)$ and develop a recursive construction to provide an infinite family of affine-invariant TQSs, that is an affine-invariant TQS $\left(p^{m}\right)$ for prime $p$ and any positive integer $m$.

Lastly, for 3 -fold quadruple systems with specific automorphism groups, Munemasa and Sawa 82 provided a perfect answer to their existence, which can be seen as a generalization of Köhler's results [67]. They proved that there exists a simple 3 -fold quadruple system with $A \rtimes \operatorname{Aut}(A)$ as its automorphism group for any abelian group $A$ of order $v \equiv 2(\bmod 4)$. Moreover, for resolvable quadruple systems, Sawa [100] considered the resolution classes under cyclic permutations, and proved that the necessary divisibility conditions for the existence of a $\lambda$-fold quadruple system with a cyclic resolution are sufficient for any $\lambda \equiv 0(\bmod 3)$.

### 1.4 Difference families and cyclotomic cosets

Let $G$ be a finite group of order $v$, written additively, and let $k$ be a positive integer. Let $A$ be a subset of $G$. Then the multiset

$$
\Delta A=\{x-y \mid x, y \in A, x \neq y\}
$$

is defined to be the list of differences of $A$.
Definition 1.4.1 (difference family). Let $\mathcal{A}$ be a collection of subsets of $G$. If every nonzero element of $G$ occurs exactly $\lambda$ times in the list $\Delta \mathcal{A}=\bigcup_{A \in \mathcal{A}} \Delta A$ then $\mathcal{A}$ is called a $(v, k, \lambda)$ difference family (DF) over $G$, where $v$ is the order of the DF. The members of $\mathcal{A}$ are called base blocks.

The number of base blocks should be $\frac{\lambda(v-1)}{k(k-1)}$, which implies the necessary divisibility condition

$$
\begin{equation*}
\lambda(v-1) \equiv 0 \quad(\bmod k(k-1)) \tag{1.4}
\end{equation*}
$$

for the existence of a $(v, k, \lambda)$-DF.
Example 1.4.2. Let

$$
B_{1}=\{0,1,3,13,34\}, B_{2}=\{0,4,9,23,45\}, \text { and } B_{3}=\{0,6,17,24,32\}
$$

Then $\left\{B_{1}, B_{2}, B_{3}\right\}$ forms a $(61,5,1)$-DF over $\mathbb{Z}_{61}$.
Moreover, if all the base blocks of a $(v, k, \lambda)$-DF, say $\mathcal{A}$, are mutually disjoint, then $\mathcal{A}$ is said to be a $(v, k, \lambda)$ disjoint difference family (DDF). DDFs have an important application, that is the construction of resolvable 2-designs via DDFs due to Ray-Chaudhuri and Wilson 98 (see also 50] Theorem 3.2.5).

A $(v, k, \lambda)$-DF over a cyclic group $\mathbb{Z}_{v}$ is simply denoted by a $(v, k, \lambda)$ cyclic difference family (CDF). By translating the base blocks of a $(v, k, \lambda)$-CDF, one can immediately obtain a strictly cyclic $2-(v, k, \lambda)$ design. Moreover, if $G$ is an elementary abelian group, that is an abelian group in which every nonzero element has the same order, then a $(v, k, \lambda)$-DF over $G$ is said to be elementary abelian. This is equivalent to considering an elementary abelian DF of order $q$ as in the additive group of the finite field $\mathbb{F}_{q}$ for a prime power $q$. The existence and
constructions of DFs over $\mathbb{Z}_{v}$ and $\mathbb{F}_{q}$ have been extensively studied as essential problems in design theory.

The existence of elementary abelian ( $q, 3,1$ )-DFs was settled by Netto 87]. For $4 \leq k \leq 6$, constructions and existence were investigated by Bose [10], Buratti [14, 15, 16, 17], and Wilson [120] over nearly six decades, and were finally settled by Chen and Zhu [28, 29]. We summarize the results as follows:

Theorem 1.4.3. For $3 \leq k \leq 6$, there exists a $(q, k, 1)$-DF over $\mathbb{F}_{q}$ for any prime power $q \equiv 1(\bmod k(k-1))$ except for $(q, k)=(61,6)$.

When $k \geq 7$, existence has not been completely determined. However, the asymptotic existence has been established by Wilson [120].

Theorem 1.4.4 (Wilson [120). Let $q$ be a prime power with $\lambda(q-1) \equiv 0$ $(\bmod k(k-1))$. Then there exists a $(q, k, \lambda)$-DF over $\mathbb{F}_{q}$ if one of the following holds:
(i) $\lambda$ is a multiple of $\frac{k}{2}$ or $\frac{k-1}{2}$,
(ii) $\lambda \geq k(k-1)$,
(iii) $q>\binom{k}{2}^{k(k-1)}$.

The bound in Theorem 1.4 .4 (iii) for the asymptotic existence of elementary abelian DFs was greatly improved by Buratti and Pasotti [21] as

$$
\begin{equation*}
q>\binom{k}{2}^{2 k} \tag{1.5}
\end{equation*}
$$

This is a consequence of Buratti and Pasotti's 21] main theorem (see Theorem 1.4.8 which can be more widely applied. We will discuss the ideas of the proofs in detail in Section 3.1.

The existence of $(v, k, 1)$-CDFs was solved for $k=3$ by Peltesohn [89] a long time ago. However, the cases when $k \geq 4$ remains unsolved so far. A recursive construction was introduced by Colbourn and Colbourn [35] and then generalized by Jimbo and Kuriki 61 utilizing the notation of cyclic difference matrices.

Definition 1.4.5 (cyclic difference matrix). A $(v, k, \lambda)$ cyclic difference matrix $(\mathrm{CDM})$ is defined to be a $k \times \lambda v$ matrix $M=\left(m_{i j}\right)$ with entries in $\mathbb{Z}_{v}$, where for any pair of indices $\left(i_{1}, i_{2}\right)$, the list of differences $\left\{m_{i_{1} j}-m_{i_{2} j} \mid 1 \leq j \leq \lambda v\right\}$ covers every element of $\mathbb{Z}_{v}$ exactly $\lambda$ times.

In particular, when $\lambda_{1}=\lambda_{2}=1$, if $u$ is an integer which is relatively prime to $(k-1)$ !, then a $(u, k, 1)$-CDM exists (see [35).

Theorem 1.4.6 (Jimbo and Kuriki [61]). If there exists a ( $\left.v, k, \lambda_{1}\right)-C D F, a$ $\left(u, k, \lambda_{1} \lambda_{2}\right)-C D F$, and $a\left(u, k, \lambda_{2}\right)-C D M$, then there exists a $\left(u v, k, \lambda_{1} \lambda_{2}\right)-C D F$.

In order to further improve the existence of DFs and "DF-like" combinatorial structures, including graph decompositions and combinatorial codes (see Lamken and Wilson [68] for a universal framework), plenty of work has been done in the past decades. In particular, for improving the asymptotic existence of DFs, Weil's Theorem on multiplicative character sums plays an essential role. Before stating the theorem, we need some notation which is used throughout this dissertation.

Let $e$ be a positive integer and $q$ be a prime power with $q \equiv 1(\bmod e)$. Let $\mathbb{F}_{q}$ and $\mathbb{F}_{q}^{*}$ denote the finite field of order $q$ and its multiplicative group, respectively. Suppose $g$ is a primitive element in $\mathbb{F}_{q}$. Let $C^{(e)}$ denote the multiplicative subgroup of $\mathbb{F}_{q}^{*}$ generated by $g^{e}$. Then, $\mathcal{C}^{(e)}:=\left\{C_{0}^{(e)}, C_{1}^{(e)}, \ldots, C_{e-1}^{(e)}\right\}$ forms a coset decomposition of $\mathbb{F}_{q}^{*}$, where $C_{i}^{(e)}:=g^{i} C^{(e)}$ is known as a cyclotomic coset of index $e$ for each $0 \leq i \leq e-1$. Moreover, let $\theta_{e}$ be a primitive $e$ th root of unity of the complex field $\mathbb{C}$ and $\chi$ be the multiplicative character of $\mathbb{F}_{q}$ of order $e$ defined by

$$
\chi(x)=\theta_{e}^{i} \text { for } x \in C_{i}^{(e)} \text { with } 0 \leq i \leq e-1 \quad \text { and } \quad \chi(0)=0
$$

Theorem 1.4.7 (Weil's Theorem, see also [75] Theorem 5.41). Let $f \in \mathbb{F}_{q}[x]$ be a polynomial that is not of the form $c g^{e}$ for some $c \in \mathbb{F}_{q}$ and $g \in \mathbb{F}_{q}[x]$. Then,

$$
\left|\sum_{x \in \mathbb{F}_{q}} \chi(f(x))\right| \leq(d-1) \sqrt{q}
$$

where $d$ is the number of distinct roots of $f$ in its splitting field over $\mathbb{F}_{q}$.
Notably, Buratti and Pasotti [21] derived an intermediate theorem from Weil's Theorem 1.4.7, which is a quite friendly and powerful tool for combinatorial studies (also given by Chang and Ji [25] independently).

Theorem 1.4.8 (Buratti and Pasotti [21] Theorem 2.2). Let $q \equiv 1(\bmod e)$ be a prime power. Let $\left\{b_{1}, b_{2}, \ldots, b_{t}\right\}$ be an arbitrary $t$-subset in $\mathbb{F}_{q}$ and let $\left(j_{1}, j_{2}, \ldots, j_{t}\right)$ be an arbitrary t-tuple of $\mathbb{Z}_{e}$. Set $X=\left\{x \in \mathbb{F}_{q} \mid x-b_{i} \in\right.$ $C_{j_{i}}^{(e)}$ for each $\left.1 \leq i \leq t\right\}$. Then, $|X|>n$ whenever $q>Q(e, t, n)$, where
$Q(e, t, n)=\left(\frac{U+\sqrt{U^{2}+4 e^{t-1}(t+e n)}}{2}\right)^{2} \quad$ and $\quad U=\sum_{h=1}^{t}\binom{t}{h}(e-1)^{h}(h-1)$.
In particular, $X$ is not empty if $q>Q(e, t):=Q(e, t, 0)$.
There is a special class of difference family which can be directly constructed from cyclotomic cosets. An elementary abelian $(q, k, 1)$ - DF is called radical if its base blocks are cosets of $C^{\left(\frac{q-1}{k}\right)}$ for odd $k$, or the union of a coset of $C^{\left(\frac{q-1}{k-1}\right)}$ and $\{0\}$ for even $k$. The terminology seems to have been introduced first by Buratti 16, but these constructions have been extensively studied by Wilson 120 in earlier times.

Theorem 1.4.9 (Wilson 120$)$. Let $q \equiv 1(\bmod k(k-1))$ be a prime power. Let $\zeta_{k}$ and $\zeta_{k-1}$ be a primitive $k$ th and $(k-1)$ th root of unity in $\mathbb{F}_{q}$, respectively, and let

$$
H= \begin{cases}\left\{\zeta_{k}^{i}-1 \left\lvert\, 1 \leq i \leq \frac{k-1}{2}\right.\right\}, & \text { if } k \text { is odd } \\ \left\{\zeta_{k-1}^{i}-1 \left\lvert\, 1 \leq i \leq \frac{k}{2}-1\right.\right\} \cup\{1\}, & \text { if } k \text { is even } .\end{cases}
$$

If $H$ forms a complete system of representatives of $\mathcal{C}^{\left(\frac{k-1}{2}\right)}$ for odd $k$ or $\mathcal{C}^{\left(\frac{k}{2}\right)}$ for even $k$, then there exists a radical $(q, k, 1)$-DF over $\mathbb{F}_{q}$.

A necessary and sufficient condition for the existence of radical $(q, k, 1)$ DFs for $k=4,5$ was established by Buratti [15], who generalized the results of Bose 10 to "perfect packing" problems (see also [17, 18). We will give a generalization of Wilson and Buratti's results on radical DFs in Section 3.2.2.

### 1.5 Grid-block designs and grid-block difference families

Let $V$ be a finite set of cardinality $v$, and $\mathcal{B}$ be a collection of $r \times k$ arrays with $r k$ distinct entries in $V$. We call the elements of $V$ and $\mathcal{B}$, respectively, points and grid-blocks. Two points are collinear in a grid-block B if they lie in the same row or in the same column of $B$.

Definition 1.5.1 (grid-block design). A pair $(V, \mathcal{B})$ is an $r \times k$ grid-block design (resp. packing, covering) on $v$ points, or a $(v, r \times k, 1)$ grid-block design (resp. packing, covering), if any pair of distinct points of $V$ is collinear in exactly (resp. at most, at least) one grid-block of $\mathcal{B}$.

In particular, when $v=r k$ and $r=k$ hold, $(V, \mathcal{B})$ is called a lattice square design. The study of lattice square designs was motivated by agricultural experiments by Yates 128 in the 1940s. Later, in 1971, Raghavarao proposed a construction of $r \times k$ grid-block designs on $p^{2}$ points when $p$ is an odd prime, in the monograph of experimental designs 97]. Hereafter, from the aspect of combinatorics, Hwang 57] proved that an $r \times k$ grid-block design on $r k$ points exists if and only if $r=k$ (i.e., a lattice square design) is odd, and $r-1$ mutually orthogonal Latin squares (MOLSs) of order $r$ exist. Moreover, Hwang 57] also proposed grid-block designs as an application to DNA library screening. Since then, Fu et al. 47] formally introduced the notion of grid-block designs with more general parameters, and discussed their further applications to life sciences as group testing models (see also 40).

Although great progress has been made in the sequencing techniques for DNA library screening, the experimental designs and data analysis on microarrays still present a higher level of statistical challenges (cf. 81] Chapter 13). Moreover, in recent years, numerous applications of group testing designs have been investigated in the areas of coding theory and computer science. For example, the connections of group testing with superimposed codes (see 41]) and compressed sensing (see 4) have attracted much attention.

Grid-block designs (resp. packings, coverings) can be naturally presented as graph decompositions. Let $H$ denote a (finite, simple, and undirected) graph. A collection $\mathcal{A}$ of subgraphs of $H$ is said to be a decomposition of $H$ if each edge of $H$ appears in exactly one subgraph in $\mathcal{A}$. Moreover, if every graph in $\mathcal{A}$ is isomorphic to a graph $G$, then $\mathcal{A}$ is said to be a $G$-decomposition of $H$. In particular, if $H=K_{v}$, namely, the complete graph on $v$ vertices, the $G$-decomposition is also known as a $G$-design of order $v$. Accordingly, under the assumption that every graph in $\mathcal{A}$ is isomorphic to $G$, if each edge of $K_{v}$ appears in at least (resp. at most) one of the graphs in $\mathcal{A}$, then $(V, \mathcal{A})$ is said to be a $G$-covering (resp. $G$-packing), where $V$ denotes the vertex set of $K_{v}$.

Let $L_{r, k}$ denote the Cartesian product graph of the complete graphs $K_{r}$ and $K_{k}$. An $r \times k$ grid-block design (resp. packing, covering) is nothing but an $L_{r, k}$-design (resp. packing, covering).

Clearly, one $2 \times 2$ grid-block design is well known as a 4 -cycle system $\left(C_{4^{-}}\right.$ design), which has been extensively studied (see [77] for details). Hence, we always suppose $\min \{r, k\} \geq 2$ and $\max \{r, k\} \geq 3$ when we discuss $r \times k$ gridblock designs. It is easy to obtain the necessary conditions for the existence of a $(v, r \times k, 1)$ grid-block design.

Proposition 1.5.2. If there exists a $(v, r \times k, 1)$ grid-block design, then

$$
\begin{equation*}
v-1 \equiv 0 \quad(\bmod r+k-2) \quad \text { and } \quad v(v-1) \equiv 0 \quad(\bmod r k(r+k-2)) . \tag{1.7}
\end{equation*}
$$

Carter 23 studied $2 \times 3$ grid-block designs as graph decompositions into $L_{2,3}$, which is considered as the only non-bipartite connected cubic graph with six vertices. As a consequence, it was proved that the necessary condition (1.7) for a $(v, 2 \times 3,1)$ grid-block design is sufficient. Moreover, for the cases of $2 \times k$, the existence problems have been completely settled for $k \in\{4,5,6\}$ by Mutoh et al. [86] $(k=4)$, Li et al. 74] $(k=5)$, and Wang and Colbourn 115 ] $(k=6)$. We summarize their results as follows:

Theorem 1.5.3. $A(v, 2 \times k, 1)$ grid-block design exists if and only if
(i) $v \equiv 1(\bmod 9)$ for $k=3$,
(ii) $v \equiv 1(\bmod 32)$ for $k=4$,
(iii) $v \equiv 1(\bmod 25)$ for $k=5$,
(iv) $v \equiv 1(\bmod 72)$ for $k=6$.

Zhang et al. [131] studied $3 \times 4$ and $4 \times 4$ grid-block designs, and proved the necessary conditions 1.7 ) are (almost) sufficient.

Theorem 1.5.4 (Zhang et al. [131]). There exists a ( $v, 4 \times k, 1$ ) grid-block design if and only if
(i) $v \equiv 1,16,21,36(\bmod 60)$ except $v=16$ and possibly except $v \in\{60 n+36 \mid$ $n=1,2,4,5,10,20,22,26\} \cup\{60 n+16 \mid n=2,3,4,7,10,18,23\}$ for $k=3$.
(ii) $v \equiv 1(\bmod 96)$ for $k=4$.

For a grid-block design (resp. packing, covering), if the set of grid-blocks can be partitioned into resolution classes, in each of which every point occurs exactly once. We will consider the constructions of resolvable grid-blocks designs (resp. packings, coverings) in Chapter 4.

Similarly to $t$-designs, we are also concerned with grid-block designs admitting prescribed permutation groups as their automorphism groups. Let $V$ be the point set. We write

$$
\mathbf{B}=\left[\begin{array}{cccc}
b_{11} & b_{12} & \cdots & b_{1 k} \\
b_{21} & b_{22} & \cdots & b_{2 k} \\
\vdots & \vdots & & \vdots \\
b_{r 1} & b_{r 2} & \cdots & b_{r k}
\end{array}\right] \quad \text { with } b_{i j} \in V
$$

to represent an $r \times k$ grid-block. When exchanging any two rows or two columns of B , the resulting grid-block is equivalent to B . Explicitly, if we regard B as a matrix, then the grid-block $\mathrm{B}^{\prime}=P \mathrm{~B} Q$ is said to be equivalent to B for any permutation matrices $P$ and $Q$. For a permutation $\sigma$ on $V$, we define $\mathrm{B}^{\sigma}=\left[b_{i j}^{\sigma}\right]_{r \times k}$. Let $(V, \mathcal{B})$ be an $r \times k$ grid-block design. If there is a grid-block equivalent to $\mathrm{B}^{\sigma}$ in $\mathcal{B}$ for any $\mathrm{B} \in \mathcal{B}$, then $(V, \mathcal{B})$ is said to be invariant under $\sigma$. Equivalently, $\sigma$ is said to be an automorphism of $(V, \mathcal{B})$. In particular, if $\sigma$ is of order $v=|V|$, then $(V, \mathcal{B})$ is cyclic. In this case, we identify $V$ with $\mathbb{Z}_{v}$ and denote $\mathrm{B}+t=\left[b_{i j}+t\right]_{r \times k}$ for any $t \in \mathbb{Z}_{v}$. Under the action of $\mathbb{Z}_{v}, \mathcal{B}$ can be partitioned into orbits. An orbit is said to be full if its length equals $v$, otherwise, short. A cyclic grid-block design containing no short orbit is said to be strictly cyclic. We can arbitrarily choose a grid-block from each orbit to represent the whole orbit. Such a representative of a (cyclic) orbit is called a (cyclic) base grid-block.

Example 1.5.5. We list the base grid-blocks of cyclic $(v, 2 \times 3,1)$ grid-block designs of small $v$ proposed in [23]. The base grid-blocks marked by $\star$ give rise to short orbits.
(i) For $v=10$, we have $\mathrm{B}_{1}=\left[\begin{array}{lll}0 & 1 & 8 \\ 6 & 5 & 3\end{array}\right]^{\star}$ as a base grid-block.
(ii) For $v=19$, we have $\mathrm{B}_{1}=\left[\begin{array}{lll}0 & 2 & 9 \\ 6 & 5 & 1\end{array}\right]$ as a base grid-block.
(iii) For $v=28$, we have $\mathrm{B}_{1}=\left[\begin{array}{ccc}0 & 1 & 6 \\ 15 & 17 & 26\end{array}\right]$ and $\mathrm{B}_{2}=\left[\begin{array}{ccc}0 & 3 & 7 \\ 14 & 21 & 17\end{array}\right]^{\star}$ as base grid-blocks.

It is difficult to construct grid-block designs when one of $r$ and $k$ is large. In the previous studies, most constructions are based on recursive ones. Small grid-block designs are utilized as "input designs" (or "ingredient designs") for

Table 1.1: Previously known cyclic $(v, r \times k, 1)$ grid-block designs

| $r \times k$ | $v$ | References |
| :--- | :--- | :--- |
| $2 \times 3$ | $10,19,28,37,46$. | Carter $[23]$ |
|  | For all $v \equiv 1(\bmod 18)$. | Wannasit and El-Zanati $[116]$ |
| $2 \times 4$ | $33,65,97,193,225,257,289,321,353$. | Mutoh et al. $[86]$ |
| $2 \times 5$ | $51,76,101$. | Li et al. $[74]$ |
| $2 \times 6$ | $73,145,433$. | Wang and Colbourn 115$]$ |
| $3 \times 3$ | For all $v \equiv 1,9(\bmod 36)$. | Fu et al. 47. |
| $3 \times 4$ | $21,61,181,421$. | Zhang et al. 131 |
| $4 \times 4$ | $97,193$. | Zhang et al. 131$]$ |

the recursions, some of which are cyclic. We list the previously known results for cyclic $(v, r \times k, 1)$ grid-block designs in Table 1.1.

Infinite families (constructions) of cyclic grid-block designs are less known. It is remarkable that the existence of cyclic $3 \times 3$ grid-block designs has been completely solved by Fu et al. [47] by giving an explicit solution. They dealt with the rows and columns of a $3 \times 3$ grid-block as base blocks of a cyclic Steiner triple system. By an ingenious arrangement of the triples, they successfully set up the direct constructions of all base grid-blocks.

Recently, a construction of cyclic $2 \times 3$ grid-block designs has been proposed by Wannasit and El-Zanati [116. As a tripartite graph having a $\rho$-tripartite labeling, $L_{2,3}$ is proved to cyclically decompose $K_{v}$ for any $v \equiv 1(\bmod 18)$. Explicitly, the base grid-blocks of a cyclic $(v, 2 \times 3,1)$ grid-design can be expressed by

$$
\mathrm{B}_{i}=\left[\begin{array}{ccc}
1 & 18 i & -18 i+8 \\
-18 i+14 & 0 & 18 i-15
\end{array}\right]
$$

where $i \in\left\{1,2, \ldots, \frac{v-1}{18}\right\}$ (see [116] Theorem 4 and Table 1).
Besides the above mentioned constructions, the method of differences is more commonly used. In another word, strictly cyclic grid-block designs can be constructed from array type analogues of difference families. For a given $r \times k$ grid-block, each row generates $k(k-1)$ differences and each column generates $r(r-1)$ differences, so a total of $r k(r+k-2)$ differences are derived. In this manner, Fu et al. 47] called the collection of base grid-blocks as a twodimensional difference family for a strictly cyclic grid-block design. Meanwhile, Mutoh, Jimbo, and Fu [85] considered their applications to resolvable grid-block designs and used the terminology "grid-block difference family". We will use the terminology "grid-block difference family" and write an $L_{r, k}$-difference family (DF), or a ( $v, L_{r, k}, 1$ )-DF for short.

Definition 1.5.6 (grid-block difference family). Let $G$ be an additive group of order $v$ and let $\mathcal{B}$ be a collection of $r \times k$ grid-blocks with entries in $G$. $\mathcal{B}$ is called an $r \times k$ grid-block difference family (DF) over $G$, or simply a $\left(v, L_{r, k}, 1\right)$-DF, if there exists exactly one pair of collinear elements in $\mathcal{B}$, say $(a, b)$, such that $a-b=x$ for any nonzero element $x$ in $G$.

It is easy to deduce the following necessary conditions for the existence of a ( $\left.v, L_{r, k}, 1\right)$-DF:

Proposition 1.5.7. If a $\left(v, L_{r, k}, 1\right)$-DF exists, then $v \equiv 1(\bmod r k(r+k-2))$.
More results on existence, construction, and characterization of $L_{r, k}$ - DFs will be proposed in Chapter 3

### 1.6 Outline of this dissertation

In the next three chapters, we will focus on the existence and construction of affine-invariant quadruple systems, grid-block difference families, and resolvable grid-block coverings.

In Chapter 2, we investigate the constructions of affine-invariant strictly cyclic Steiner quadruple systems (AsSQSs) and affine-invariant 2-fold quadruple systems (TQSs). For a prime $p \equiv 1(\bmod 4)$, Direct Construction A establishes an $\operatorname{AsSQS}(2 p)$, provided that a 1-factor of a graph exists, where the graph is defined by using a system of generators of the projective special linear group $\operatorname{PSL}(2, p)$. Direct Construction B gives an $\operatorname{AsSQS}(2 p)$ which is 2-chromatic, provided that a rainbow 1-factor of a specific hypergraph exists. Accordingly, by proposing two recursive constructions of an $\operatorname{AsSQSs}\left(2 p^{m}\right)$ for a positive integer $m$, we prove that an $\operatorname{AsSQS}\left(2 p^{m}\right)$ exists, if the criteria developed for an $\operatorname{AsSQS}(2 p)$ are satisfied. In a similar way, the direct construction and recursive construction for affine-invariant TQSs are also given.

Chapters 3 and 4 are devoted to considering grid-block designs with the cyclic property and resolvability. In Chapter 3, we concentrate on grid-block difference families (DFs), which can be viewed as two-dimensional generalizations of DFs. Firstly, we give an intermediate algebraic consequence on the existence bound of an element satisfying certain cyclotomic conditions in a finite field. In many cases, this approach improves the bound due to Buratti and Pasotti [21]. In particular, this approach will be applied to improving the existence bound for grid-block DFs. Secondly, by considering Kronecker density via algebraic number theory, a series of cyclotomic constructions and characterizations of row-radical grid-block DFs are presented.

In Chapter 4, a construction of resolvable grid-block designs, packings, or coverings via grid-block DFs is proposed. Moreover, the optimality of grid-block covering is discussed and the optimal construction of $2 \times 3$ grid-block covering is provided.

In Chapter 5, we give concluding remarks and some problems for future study.

## Chapter 2

## Affine-invariant quadruple systems

In this chapter, we focus mainly on affine-invariant SQSs and TQSs, that is, cyclic 3- $(v, 4, \lambda)$ designs with $\lambda \in\{1,2\}$ over $\mathbb{Z}_{v}$ which admit every unit of $\mathbb{Z}_{v}$ as a multiplier. We give two direct constructions for an $\operatorname{AsSQS}(2 p)$, and two recursive constructions for an $\operatorname{AsSQS}\left(2 p^{m}\right)$, where $p \equiv 1,5(\bmod 12)$ is prime and $m$ is a positive integer.

In Section 2.1. we introduce two families of graphs, namely, "LG graphs" and "CG graphs". An LG graph is defined on a 1-dimensional projective line (which is abbreviated to the letter "L") over a finite field $\mathbb{F}_{q}$. A CG graph is defined on the cross-ratio classes (which is abbreviated to the letter "C") of a projective line. The adjacencies in both LG and CG graphs are established by a set of generators of the projective special linear group $\operatorname{PSL}(2, q)$. This new perspective would provide a possible way for making use of geometric group theory or combinatorial group theory to attack the complete proof of the existence of sSQSs and AsSQSs. In Section 2.1.3, we describe the relation between our CG graphs and "Köhler orbit graphs".

In Section 2.2.1. we give a presentation of blocks (quadruples and triples) over $\mathbb{Z}_{\frac{v}{2}} \times \mathbb{Z}_{2}$ to simplify our constructions. Then, under the above presentation, we present two direct Constructions 2.2.6 and 2.2.20 in Sections 2.2.2 and 2.2 .3 , respectively, where the former requires 1 -factors of CG graphs defined in Section 2.1, and the latter is related to a hypergraph which can be regarded as a pairwise balanced design (PBD). In addition, the sSQSs obtained from Construction 2.2 .20 are 2-chromatic, so that a few unknown parameters of 2 -chromatic SQSs can be determined.

We use Section 2.3 .1 to summarize some notation and useful preliminaries for the constructions below. Two recursive constructions are presented in Sections 2.3.2 and 2.3.3 showing that an $\operatorname{AsSQS}\left(2 p^{m}\right)$ can be constructed via an AsSQS $(2 p)$ derived from direct Constructions 2.2 .6 and 2.2 .20

In Section 2.4, we prove a necessary condition for the existence of an $\operatorname{AsSQS}(v)$
for $v \equiv 2,10(\bmod 24)$, and thereby establish a non-existence result.
Furthermore, direct and recursive constructions are presented for affineinvariant TQSs in Sections 2.5 and 2.6, respectively.

It is mentioned that sSQSs are closely related to optical orthogonal codes (OOCs). In fact, combinatorial designs, in particular cyclic designs, are widely used in many other areas, for example, designs of experiments, group testing [39], authentication codes [88, 110], filing schemes [6, 126], etc. For the applications to these areas, it is desirable to generate the blocks with less storage and time. Also, for a given $t$-subset (for example, a pair, a triple, etc.) $T$ of the point set $V$, it is usually required to find the blocks containing the certain $T$. The affine-invariant property works effectively for these problems. Finally, Section 2.7 is devoted to giving a brief explanation on these approaches.

Throughout this chapter, we always suppose $p$ is a prime satisfying $p \equiv 1$ (mod 4). Besides the standard notation, we use the symbols $\uplus$ and $\cup$ to denote the union of multisets and the disjoint union of sets, respectively.

### 2.1 Graphs associated with $\operatorname{PSL}(2, q)$

In this section, we define two families of graphs, namely, LG graphs and CG graph. An LG graph can be defined on any finite set $V$ with a group acting on it. A CG graph can be derived from a specific LG graph defined on the projective line over a finite field, and plays an essential role for our constructions for affine-invariant sSQSs and TQSs.

Hereafter, when saying a graph, we mean an undirected graph in which multiple edges and self-loops are allowed. More precisely, the edge set of a graph is considered as a multi-set, and a self-loop is represented by a singleton in the edge set. The degree of a vertex, say $x$, is defined as the number of edges, including multiple edges and self-loops, which contain $x$.

For more notion of graphs and hypergraphs, the reader is referred to the textbooks [26, 37] for details.

### 2.1.1 LG graphs

First, we introduce the most general definition of an LG graph and present a series of its basic properties.

Definition 2.1.1 (LG graph). Let $V$ be a finite set, and let $G$ be a group acting on $V$ such that $V^{G} \subseteq V$. Let $\Sigma$ be a finite subset of $G$ consisting of involutions. $\mathrm{LG}(V, \Sigma)$ is defined to be the graph $(V, E)$ with edge set $E=\left\{\left\{x, x^{\sigma}\right\} \mid x \in\right.$ $V, \sigma \in \Sigma\}$. Multiple edges and self-loops are allowed, and thus $E$ is treated as a multiset.

Proposition 2.1.2. With the notation in Definition 2.1.1, the following hold:
(i) If $\Sigma$ consists of involutions in $G$ and $\Sigma^{\prime} \subset \Sigma$, then $\mathrm{LG}\left(V, \Sigma^{\prime}\right)$ is an edgeinduced subgraph of $\mathrm{LG}(V, \Sigma)$.
(ii) $\mathrm{LG}(V, \Sigma)$ has a self-loop at vertex $x \in V$ if and only if $x$ is a fixed point of some $\sigma \in \Sigma$.
(iii) Suppose $|\Sigma| \geq 2$. $\operatorname{LG}(V, \Sigma)$ has multiple edges $\{x, y\} \in E$ only if both $x$ and $y$ are fixed points of $\sigma_{1} \sigma_{2}$ for distinct $\sigma_{1}, \sigma_{2} \in \Sigma$, where the group $G$ is written multiplicatively.
(iv) Every vertex of $\operatorname{LG}(V, \Sigma)$ is of degree $|\Sigma|$.
(v) $\mathrm{LG}(V, \Sigma)$ consists of $r$ vertex-disjoint subgraphs, each of which has a vertex set identical with an orbit of $V$ under the action of $\langle\Sigma\rangle$, where $r$ is the number of orbits.
(vi) If $G$ acts transitively on $V$ and $\Sigma$ is a generating set of $G$, then $\operatorname{LG}(V, \Sigma)$ is connected.

Proof. (i) and (ii) follow straightforwardly from Definition 2.1.1.
(iii) By Definition 2.1.1, $\{x, y\} \in E$ if and only if $y=x^{\sigma}$ for some $\sigma \in \Sigma$. Suppose $\{x, y\}$ appears more than once in $E$. In another word, there exist $\sigma_{1}, \sigma_{2} \in \Sigma$, such that $x^{\sigma_{1}}=x^{\sigma_{2}}=y$. Since $\sigma_{1}$ and $\sigma_{2}$ are involutions, we also have $y^{\sigma_{1}}=y^{\sigma_{2}}=x$. This implies both $x$ and $y$ are fixed points of $\sigma_{1} \sigma_{2}$.
(iv) In fact, any edge of $\operatorname{LG}(V, \Sigma)$ is an orbit of $V$ under the action of some $\sigma \in \Sigma$. Hence, for any $\sigma \in \Sigma$, each vertex has degree 1 in the subgraph $\mathrm{LG}(V,\{\sigma\})$ (cf. (ii)). When multiple edges and self-loops are counted (cf. (iii) and (iii)), it is clear that each vertex of $\operatorname{LG}(V, \Sigma)$ has degree $|\Sigma|$.
(v) Let $V_{1}, V_{2}, \ldots, V_{r}$ denote the orbits of $V$ under the action of $\langle\Sigma\rangle$. Then, $\mathrm{LG}\left(V_{i}, \Sigma\right)$ is obviously a vertex-induced subgraph of $\mathrm{LG}(V, \Sigma)$ for each $1 \leq i \leq r$. Let $E_{i}$ denote the edge set of $\operatorname{LG}\left(V_{i}, \Sigma\right)$. Note that $x$ and $x^{\sigma}$ always lie in the same orbit under the action of $\langle\Sigma\rangle$ for any $x \in V$ and $\sigma \in \Sigma$. Thus, we have $E=\left\{\left\{x, x^{\sigma}\right\} \mid x \in \biguplus_{i=1}^{r} V_{i}, \sigma \in \Sigma\right\}=$ $\biguplus_{i=1}^{r}\left\{\left\{x, x^{\sigma}\right\} \mid x \in V_{i}, \sigma \in \Sigma\right\}=\biguplus_{i=1}^{r} E_{i}$.
(vi) Every $\tau \in G$ can be represented by a sequence of generators in $\Sigma$, say $\tau=\sigma_{\ell} \sigma_{\ell-1} \cdots \sigma_{1}$, where $\sigma_{i} \in \Sigma$ for $1 \leq i \leq \ell$. In terms of graphs, there exists a walk from $u$ to $u^{\tau}$ of length $\ell$ for any vertex $u$ of $\operatorname{LG}(V, \Sigma)$. Conversely, by the transitivity of $G$, for any $u, v \in V$, there must exist $\tau \in G$, such that $v=u^{\tau}$. In other words, there exists a walk from $u$ to $v$ for any distinct vertices $u, v$. Therefore, $\operatorname{LG}(V, \Sigma)$ is connected.

Let $\mathbf{P}^{1}\left(\mathbb{F}_{q}\right)$ denote the projective line over $\mathbb{F}_{q}$ which can be identified with the line $\mathbb{F}_{q}$ extended by a point at infinity, namely, $\mathbb{F}_{q} \cup\{\infty\}$. Let

$$
\sigma: x \mapsto \frac{a x+b}{c x+d}
$$

be a fractional linear transformation on $\mathbf{P}^{1}\left(\mathbb{F}_{q}\right)$ with $a, b, c, d \in \mathbb{F}_{q}$ and $a d-b c=$ 1. Then all such transformations form a group under composition which is known as the projective special linear group of degree 2 , and is denoted by $\operatorname{PSL}(2, q)$, where $\sigma(\infty)=\frac{a}{c}, \sigma\left(-\frac{d}{c}\right)=\infty$ if $c \neq 0$, and $\sigma(\infty)=\infty$ if $c=0$.

In particular, take

$$
\begin{equation*}
\sigma_{\boldsymbol{A}}: x \mapsto 1-x, \quad \sigma_{B}: x \mapsto \frac{1}{x}, \quad \text { and } \quad \sigma_{C}: x \mapsto \frac{1-x}{1-2 x} \tag{2.1}
\end{equation*}
$$

It is easy to see that $\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}$, and $\sigma_{\boldsymbol{C}}$ are involutions. Now we begin to consider the graph $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{q}\right),\left\{\sigma_{A}, \sigma_{B}, \sigma_{C}\right\}\right)$.

Example 2.1.3. The graphs $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$ for $p=13$ and $p=29$ are shown in Figures 2.1 and 2.2 , respectively, where the edges labeled by $\boldsymbol{A}, \boldsymbol{B}$, and $\boldsymbol{C}$ are contained in $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}\right\}\right), \operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{B}}\right\}\right)$, and $\mathrm{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{C}\right\}\right)$, respectively.


Figure 2.1: The graph $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{13}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$


Figure 2.2: The graph $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{29}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$

Example 2.1.4. The graph $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{q}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$ for $q=25$ is shown in Figures 2.3 , where $\mathbb{F}_{25}$ is considered as $\mathbb{F}_{5}[\alpha] /\left(\alpha^{2}+2\right)$, and the edges labeled by $\boldsymbol{A}, \boldsymbol{B}$, and $\boldsymbol{C}$ are contained in $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{q}\right),\left\{\sigma_{\boldsymbol{A}}\right\}\right), \operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{q}\right),\left\{\sigma_{\boldsymbol{B}}\right\}\right)$, and $\mathrm{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{q}\right),\left\{\sigma_{C}\right\}\right)$, respectively. This graph consists of two connected components, one of which is $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{5}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$.


Figure 2.3: The graph $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{25}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$
Lemma 2.1.5. Let $p \equiv 1(\bmod 4)$ be a prime. Then $\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}$ is a generating set of $\operatorname{PSL}(2, p)$, where $\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}$, and $\sigma_{\boldsymbol{C}}$ follow the definitions in 2.1.

Proof. The group $\operatorname{PSL}(2, p)$ can also be interpreted as a matrix group, namely,

$$
\operatorname{PSL}(2, p)=\left\{\left.M=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \right\rvert\, a, b, c, d \in \mathbb{F}_{p}, \operatorname{det} M=1\right\}
$$

For any $\delta \in \mathbb{F}_{p}^{*}$, since $\frac{a x+b}{c x+d}=\frac{\delta a x+\delta b}{\delta c x+\delta d}$, we identify $\delta M$ with $M$ when $M \in$ $\operatorname{PSL}(2, p)$ is considered as a matrix. In this proof, we use the matrix representation. Then $\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}$, and $\sigma_{\boldsymbol{C}}$ correspond to three matrices over $\mathbb{F}_{p}$ respectively, namely,

$$
\boldsymbol{A}=\left(\begin{array}{cc}
-1 & 1 \\
0 & 1
\end{array}\right), \quad \boldsymbol{B}=\left(\begin{array}{cc}
0 & 1 \\
1 & 0
\end{array}\right), \quad \text { and } \quad \boldsymbol{C}=\left(\begin{array}{cc}
1 & -1 \\
2 & -1
\end{array}\right)
$$

Note that -1 is a square in $\mathbb{F}_{p}$ if and only if $p \equiv 1(\bmod 4)$. Since $\operatorname{det} \boldsymbol{A}=$ $\operatorname{det} \boldsymbol{B}=-1$ and $\operatorname{det} \boldsymbol{C}=1$, by identifying $\boldsymbol{A}$ (resp. $\boldsymbol{B}$ ) with $\sqrt{-1} \boldsymbol{A}$ (resp. $\sqrt{-1} \boldsymbol{B})$, we have $\boldsymbol{A}, \boldsymbol{B}, \boldsymbol{C} \in \operatorname{PSL}(2, p)$. To complete the proof, we employ a theorem due to Behr and Mennicke [5] (see also Coxeter and Moser [36] § 7.5)
who showed that, for odd prime $p, \operatorname{PSL}(2, p)$ can be represented by the system of generators and relations

$$
\begin{equation*}
\boldsymbol{S}^{p}=\boldsymbol{T}^{2}=(\boldsymbol{S T})^{3}=\left(\boldsymbol{S}^{2} \boldsymbol{T} \boldsymbol{S}^{\frac{1}{2}(p+1)} \boldsymbol{T}\right)^{3}=\boldsymbol{I} \tag{2.2}
\end{equation*}
$$

where $\boldsymbol{I}$ is the identity matrix. Let $\boldsymbol{S}=\boldsymbol{C B} \boldsymbol{A}=\left(\begin{array}{ll}1 & 0 \\ 1 & 1\end{array}\right)$ and $\boldsymbol{T}=\boldsymbol{A} \boldsymbol{B} \boldsymbol{C} \boldsymbol{B} \boldsymbol{A}=$ $\left(\begin{array}{cc}0 & 1 \\ -1 & 0\end{array}\right)$. Since $\langle\boldsymbol{S}, \boldsymbol{T}\rangle \leq\langle\boldsymbol{A}, \boldsymbol{B}, \boldsymbol{C}\rangle \leq \mathrm{PSL}(2, p)$, it remains to show that $\boldsymbol{S}, \boldsymbol{T}$ satisfy $(2.2)$. It is easy to see $\boldsymbol{T}^{2}=-\boldsymbol{I}, \boldsymbol{S}^{p} \equiv \boldsymbol{I}(\bmod p)$, and $\left(\boldsymbol{S}^{2} \boldsymbol{T} \boldsymbol{S}^{\frac{1}{2}(p+1)} \boldsymbol{T}\right)^{3} \equiv$ $\boldsymbol{I}(\bmod p)$, where $\boldsymbol{I}$ and $-\boldsymbol{I}$ are identical in $\operatorname{PSL}(2, p)$. Thus $\langle\boldsymbol{S}, \boldsymbol{T}\rangle=\operatorname{PSL}(2, p)$ which implies $\langle\boldsymbol{A}, \boldsymbol{B}, \boldsymbol{C}\rangle=\operatorname{PSL}(2, p)$.

In order to examine the special vertices of $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{q}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$, we need the following lemma:

Lemma 2.1.6. Let $q$ be an odd prime power. Then,
(i) There exists $\chi \in \mathbb{F}_{q}$ such that $\chi^{\sigma_{C}}=\chi$ if and only if $q \equiv 1(\bmod 4)$.
(ii) There exists $\mu \in \mathbb{F}_{q}$ such that $\mu^{\sigma_{C}}=\mu^{\sigma_{B}}$ if and only if $q \equiv \pm 1(\bmod 5)$.
(iii) There exists $\xi \in \mathbb{F}_{q}$ such that $\xi^{\sigma_{A}}=\xi^{\sigma_{B}}$ if and only if $q \equiv 1(\bmod 3)$.

Proof. This is equivalent to studying the roots of the equations

$$
\begin{equation*}
2 \chi^{2}-2 \chi+1=0, \quad \mu^{2}-3 \mu+1=0, \quad \text { and } \quad \xi^{2}-\xi+1=0 \tag{2.3}
\end{equation*}
$$

which can be formally written by

$$
\begin{equation*}
\chi=\frac{1+\sqrt{-1}}{2}, \quad \mu=\frac{3+\sqrt{5}}{2}, \quad \text { and } \quad \xi=\frac{1+\sqrt{-3}}{2} . \tag{2.4}
\end{equation*}
$$

These expressions respectively require $-1,5$, and -3 to be squares in $\mathbb{F}_{q}$. If $q=p$ is a prime, it can be obtained by using the law of quadratic reciprocity that -1 (resp., 5 or -3 ) is a square modulo $p$ if and only if $p \equiv 1(\bmod 4)$ $($ resp., $p \equiv \pm 1(\bmod 5)$ or $p \equiv 1(\bmod 3))$. In the extension field $\mathbb{F}_{q}$ with $q=p^{n}, n \geq 1$, if $p \equiv 1(\bmod 4)$, then -1 is known to be a square in the subfield $\mathbb{F}_{p} \subset \mathbb{F}_{q}$. On the other hand, if $p \equiv 3(\bmod 4)$, then the polynomial $x^{2}+1$ is irreducible over $\mathbb{F}_{p}$, and its splitting field is nothing but $\mathbb{F}_{p^{2}}$. Hence, -1 is a square in $\mathbb{F}_{p^{2}}$ when $p \equiv-1(\bmod 4)$. In summary, -1 is a square in $\mathbb{F}_{p^{n}}$ with $p$ odd, if and only if $p \equiv 1(\bmod 4)$ or $n$ is even, that is, $p^{n} \equiv 1(\bmod 4)$. In the same manner, it can be shown that 5 and -3 are squares in $\mathbb{F}_{q}$ if and only if $q \equiv \pm 1(\bmod 5)$ and $q \equiv 1(\bmod 3)$, respectively.

As a direct consequence of Proposition 2.1.2, we summarize the properties of $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$ as follows:

Proposition 2.1.7. The following properties of $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$ hold for $p \equiv 1(\bmod 4)$.
(i) There is a self-loop at vertex $x \in V$ if and only if $x \in\left\{2^{-1}, \infty\right\} \cup\{1,-1\} \cup$ $\{\chi, 1-\chi\}$, where $\chi=\frac{1+\sqrt{-1}}{2}$ is a root of $2 \chi^{2}-2 \chi+1=0$.
(ii) There are double edges $\{0,1\}$. If $p \equiv 1(\bmod 3)$, there are double edges $\{\xi, 1-\xi\}$, where $\xi=\frac{1+\sqrt{-3}}{2}$ is a root of $\xi^{2}-\xi+1=0$. If $p \equiv \pm 1$ $(\bmod 5)$, there are double edges $\left\{\mu, \mu^{-1}\right\}$, where $\mu=\frac{3+\sqrt{5}}{2}$ is a root of $\mu^{2}-3 \mu+1=0$.
(iii) Every vertex is of degree 3.
(iv) $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$ is connected.

Proof. (ii) and (ii) follow from Proposition 2.1 .2 (iii) and (iii) by noting that the fixed points of $\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}, \sigma_{\boldsymbol{A}} \sigma_{\boldsymbol{C}}, \sigma_{\boldsymbol{A}} \sigma_{\boldsymbol{B}}$, and $\sigma_{\boldsymbol{B}} \sigma_{\boldsymbol{C}}$ are $\left\{2^{-1}, \infty\right\},\{1,-1\}$, $\{\chi, 1-\chi\},\{0,1\},\{\xi, 1-\xi\}$, and $\left\{\mu, \mu^{-1}\right\}$, respectively, where the existence criteria for $\chi, \mu$, and $\xi$ agree with Lemma 2.1.6.
(iii) is a direct conclusion of Proposition 2.1.2 (iv). Lemma 2.1.5 indicates that $\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}$ is a generating set of $\operatorname{PSL}(2, p)$. Moreover, $\operatorname{PSL}(2, p)$ is wellknown to be transitive. Thus, $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{C}\right\}\right)$ is connected by Proposition 2.1.2 (vi).

### 2.1.2 CG graphs

Let $x \in \mathbf{P}^{1}\left(\mathbb{F}_{q}\right)$ and let $C(x)$ denote the orbit of $x$ under the action of the subgroups $\left\langle\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}\right\rangle$, i.e.,

$$
\begin{equation*}
C(x)=\left\{x^{\sigma} \mid \sigma \in\left\langle\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}\right\rangle\right\}=\left\{x, \frac{1}{x}, \frac{x-1}{x}, \frac{x}{x-1}, \frac{1}{1-x}, 1-x\right\} \tag{2.5}
\end{equation*}
$$

In projective geometry, $C(x)$ is also known as the cross-ratio class with respect to $x$. Since $C(x)$ is essentially a $\left\langle\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}\right\rangle$-orbit of $x \in \mathbf{P}^{1}\left(\mathbb{F}_{q}\right)$, we have

$$
|C(x)|= \begin{cases}2 & \text { if } x \in\{\xi, 1-\xi\}  \tag{2.6}\\ 3 & \text { if } x \in\{0,1, \infty\} \cup\left\{-1,2,2^{-1}\right\} \\ 6 & \text { otherwise }\end{cases}
$$

For $p \equiv 5(\bmod 12), \mathbf{P}^{1}\left(\mathbb{F}_{q}\right)$ can be partitioned into $\{0,1, \infty\},\left\{-1,2,2^{-1}\right\}$, and other $\frac{q-5}{6}$ cross-ratio classes of size 6 . For $p \equiv 1(\bmod 12), \mathbf{P}^{1}\left(\mathbb{F}_{q}\right)$ can be partitioned into $\{0,1, \infty\},\left\{-1,2,2^{-1}\right\},\{\xi, 1-\xi\}$, and other $\frac{q-7}{6}$ cross-ratio classes of size 6 , where $\xi=\frac{1+\sqrt{-3}}{2}$ is a root of $\xi^{2}-\xi+1=0$.

We remove the cross-ratio classes of odd sizes and let

$$
\begin{equation*}
\Omega_{q}=\mathbb{F}_{q} \backslash\left\{0,1,-1,2,2^{-1}\right\} \tag{2.7}
\end{equation*}
$$

We need a lemma to ensure Definition 2.1.9 is well-defined.

Lemma 2.1.8. For any $x \in \Omega_{q}$, let $R(x)=\left\{\{u, v\} \mid v=u^{\sigma_{C}}, u \in C(x), v \in\right.$ $\left.C\left(x^{\sigma_{C}}\right)\right\}$. Then the cardinality of $R(x)$ is either 2 or 4 . In particular, if $C(x)=$ $C\left(x^{\sigma_{C}}\right)$, then $|R(x)|=2$ and $x \in\left\{\mu, \mu^{-1}, 1-\mu, 1-\mu^{-1}, \chi, 1-\chi\right\}$, where $\mu=\frac{3+\sqrt{5}}{2}$ is a root of $\mu^{2}-3 \mu+1=0$, and $\chi=\frac{1+\sqrt{-1}}{2}$ is a root of $2 \chi^{2}-2 \chi+1=0$.

Proof. It is clear that $\left\{x, x^{\sigma_{C}}\right\} \in R(x)$. Moreover, we have $\left\{1-x, 1-x^{\sigma_{C}}\right\} \in$ $R(x)$. Hence, $|R(x)|$ is even unless $\left\{x, x^{\sigma_{C}}\right\}=1-\left\{x, x^{\sigma_{C}}\right\}$, which implies $x \in\left\{0,1,2^{-1}\right\}$. Therefore, $|R(x)| \geq 2$ is even for any $x \in \Omega_{q}$.

If $C(x)$ coincides with $C\left(x^{\sigma_{C}}\right)$, then there must exist $\tau \in\left\langle\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}\right\rangle$ such that $x^{\tau}=x^{\sigma_{C}}$, which implies $x \in\left\{\mu, \mu^{-1}, 1-\mu, 1-\mu^{-1}, 0,1, \chi, 1-\chi\right\}$. Note that $x$ is in $\Omega_{q}$ which does not contain $\{0,1\}$. Thus, for any $x \in\left\{\mu, \mu^{-1}, 1-\mu, 1-\mu^{-1}\right\}$, $R(x)=\left\{\left\{\mu, \mu^{-1}\right\},\left\{1-\mu, 1-\mu^{-1}\right\}\right\}$. For any $x \in\{\chi, 1-\chi\}, R(x)=\{\{\chi\},\{1-$ $\chi\}\}$. In both cases, we have $|R(x)|=2$.

Then, suppose $C(x)$ is disjoint from $C\left(x^{\sigma_{C}}\right)$ and $R(x) \geq 4$. It should satisfy $\left(x^{\sigma_{B}}\right)^{\sigma_{C}}=\left(x^{\sigma_{C}}\right)^{\tau}$ or $\left((1-x)^{\sigma_{B}}\right)^{\sigma_{C}}=\left(x^{\sigma_{C}}\right)^{\tau}$ for some $\tau \in\left\langle\sigma_{\boldsymbol{A}}, \sigma_{B}\right\rangle$. For $x \in \Omega_{q}$, this criterion is satisfied only if $x \in C\left(\frac{1}{\sqrt{2}}\right) \cup C\left(-\frac{1}{\sqrt{2}}\right)$ when $q \equiv \pm 1$ $(\bmod 8)$. In this case, we have $|R(x)|=4$.

Definition 2.1.9 (CG graph). Let $X$ be a subset of $\Omega_{q}$. Let $V=\{C(x) \mid x \in$ $X\}$ and let $E$ be a multiset of subsets of $V$ consisting of

$$
\left\{C(x), C\left(x^{\sigma_{C}}\right)\right\} \text { for any } x \in X
$$

with multiplicity $\frac{1}{2} \#\left\{\{u, v\} \mid v=u^{\sigma_{C}}, u \in C(x), v \in C\left(x^{\sigma_{C}}\right)\right\}$. Then $(V, E)$ is an incidence structure which can be seen as a graph with multiple edges and self-loops, and is denoted by $\operatorname{CG}(X)$.

In particular, using $\operatorname{LG}\left(\Omega_{q},\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$, we have an equivalent definition of $\mathrm{CG}\left(\Omega_{q}\right)$. First, contract each component (corresponding to a cross-ratio class) of $\operatorname{LG}\left(\Omega_{q},\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}\right\}\right)$ into a single vertex. Then, contract each pair of edges in $\mathrm{LG}\left(\Omega_{q},\left\{\sigma_{C}\right\}\right)$ into a single edge if they are contained in the same component of $\operatorname{LG}\left(\Omega_{q},\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{C}}\right\}\right)$. The resulting graph is $\mathrm{CG}\left(\Omega_{q}\right)$.

When $p \equiv 1(\bmod 4)$, the connectivity of $\operatorname{CG}\left(\Omega_{p}\right)$ can be easily derived from the connectivity of the vertex-induced subgraph of $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$ on $\Omega_{p}$, where $\operatorname{LG}\left(\mathbf{P}^{1}\left(\mathbb{F}_{p}\right),\left\{\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{B}}, \sigma_{\boldsymbol{C}}\right\}\right)$ is connected by Proposition 2.1.7 (iv).

Since $\Omega_{p}=\emptyset$ when $p=5$, we assume $p>5$ for the following discussion of $\mathrm{CG}\left(\Omega_{p}\right)$.

Proposition 2.1.10. For prime $p \equiv 1(\bmod 4)$ and $p>5$, the following properties on $\operatorname{CG}\left(\Omega_{p}\right)$ hold:
(i) There is a self-loop at vertex $C$ if and only if $C=C(\chi)$ or $C=C(\mu)$, where $\chi=\frac{1+\sqrt{-1}}{2}$ is a root of $2 \chi^{2}-2 \chi+1=0$, and $\mu=\frac{3+\sqrt{5}}{2}$ is a root of $\mu^{2}-3 \mu+1=0$ which requires $p \equiv 1,29,41,49(\bmod 60)$.
(ii) Every vertex has degree 3 except $C(3), C(\mu)$, and $C(\xi)$, where $C(3)$ and $C(\mu)$ are of degree 2 , and $C(\xi)$ is of degree 1 , where where $\xi=\frac{1+\sqrt{-3}}{2}$ is a root of $\xi^{2}-\xi+1=0$;

Table 2.1: Special vertices of $\mathrm{CG}\left(\Omega_{p}\right)$ for $p \equiv 1(\bmod 4)$

|  | $C(3)$ | $C(\chi)$ | $C(\mu)$ | $C(\xi)$ |
| ---: | :---: | :---: | :---: | :---: |
| $p \equiv 29,41(\bmod 60)$ | $\checkmark$ | $\checkmark$ | $\checkmark$ |  |
| $p \equiv 17,53(\bmod 60)$ | $\checkmark$ | $\checkmark$ |  | $\checkmark$ |
| $p \equiv 1,49(\bmod 60)$ | $\checkmark$ | $\checkmark$ | $\checkmark$ | $\checkmark$ |
| $p \equiv 13,37(\bmod 60)$ | $\checkmark$ | $\checkmark$ |  | 1 |
| Degree | 2 | 3 | 2 | $\checkmark$ |
| Incident with a self-loop |  | $\checkmark$ | $\checkmark$ |  |
| Remarks (cf. Lemma 2.1 .6$)$ | $3^{\sigma_{B}}=2^{\sigma_{C}}$ | $\chi^{\sigma_{C}}=\chi$ | $\mu^{\sigma_{C}}=\mu^{\sigma_{B}}$ | $\xi^{\sigma_{A}}=\xi^{\sigma_{B}}$ |

Proof. (i) $\{C(x)\}$ forms a self-loop if and only if there exists $u \in C(x)$, such that $u^{\sigma_{C}} \in C(x)$, which implies $C(x)=C(\chi)$ or $C(\mu)$ by Lemma 2.1.8, In addition, 5 is a square modulo $p$ if and only if $p \equiv \pm 1(\bmod 5)$. Combined with $p \equiv 1,5(\bmod 12)$, we have $p \equiv 1,29,41,49(\bmod 60)$.
(ii) Let $\operatorname{deg} C(x)$ denote the degree of $C(x)$ in $\mathrm{CG}\left(\Omega_{q}\right)$. It is easy to see that $\operatorname{deg} C(\mu)=2$, and $\operatorname{deg} C(\chi)=3$ except when $q=13$, in which case $C(\chi)=C(3)$ is of degree 2 . Now we suppose $C(x)$ is not incident with a self-loop. By Definition 2.1.9 of $\mathrm{CG}\left(\Omega_{q}\right)$, we have

$$
\begin{align*}
\operatorname{deg} C(x) & =\frac{1}{2} \#\left\{C\left(u^{\sigma_{C}}\right) \text { is a vertex other than } C(x) \mid u \in C(x)\right\} \\
& =\frac{1}{2}|C(x)|-\frac{1}{2} \#\left\{u^{\sigma_{C}} \notin \Omega_{q} \mid u \in C(x)\right\} . \tag{2.8}
\end{align*}
$$

For any $x \in \Omega_{q}, x^{\sigma_{C}} \notin \Omega_{q}$ if and only if $x \in\left\{3^{-1}, 1-3^{-1}\right\} \subset C(3)$. Hence, it follows from (2.6) and 2.8) that $\operatorname{deg} C(3)=2, \operatorname{deg} C(\xi)=1$, and $\operatorname{deg} C(x)=3$ for any other $C(x)$ without a self-loop.

We summarize the degrees of $\operatorname{CG}\left(\Omega_{p}\right)$ in Table 2.1. For showing the existence of 1-factors of $\mathrm{CG}\left(\Omega_{p}\right)$, we need the following theorem, which can be seen as a stronger version of the famous Petersen's theorem [90] in graph theory:

Theorem 2.1.11 (Plesník 94, see also [130] Theorem 1.44, [1] Theorem 2.39). Any 2-edge-connected 3-regular graph (multigraphs without self-loops) has a 1factor excluding any pair of edges.

Theorem 2.1.12. For $p \equiv 1,5(\bmod 12)$ and $p \not \equiv 1,49(\bmod 60), \mathrm{CG}\left(\Omega_{p}\right)$ has a 1-factor if there is no bridge except its pendant edges.

Proof. Removing the self-loops and adding some auxiliary edges in $\operatorname{CG}\left(\Omega_{p}\right)$, we can obtain a 3-regular graph as follows:
(a) If $p \equiv 29,41(\bmod 60)$, we add auxiliary edges $\{C(3), C(\mu)\}$ and $\{C(\chi), C(\mu)\}$.
(b) If $p \equiv 5(\bmod 12)$ and $p \not \equiv 29,41(\bmod 60)$, we add an auxiliary edge $\{C(3), C(\chi)\}$.


Figure 2.4: $\mathrm{CG}\left(\Omega_{13}\right)$ and $\mathrm{CG}\left(\Omega_{17}\right)$


Figure 2.5: $\operatorname{CG}\left(\Omega_{29}\right), \operatorname{CG}\left(\Omega_{37}\right)$, and $\operatorname{CG}\left(\Omega_{41}\right)$
(c) If $p \equiv 1(\bmod 12)$ and $p \not \equiv 1,49(\bmod 60)$, we add auxiliary edges $\{C(3), C(\xi)\}$ and $\{C(\chi), C(\xi)\}$.

We denote the resulting graph by $\mathrm{CG}^{\dagger}\left(\Omega_{p}\right)$. Since there are no more than two auxiliary edges for all these cases, by Theorem 2.1.11, if $\mathrm{CG}^{\dagger}\left(\Omega_{p}\right)$ is bridgeless (2-edge-connected), there must exist a 1 -factor excluding all auxiliary edges, which is a 1-factor of $\operatorname{CG}\left(\Omega_{p}\right)$ as well.

The following results are verified by computers:
Lemma 2.1.13. $\mathrm{CG}\left(\Omega_{p}\right)$ has a 1 -factor for all primes $p<10^{5}$ with $p \equiv 1$ $(\bmod 4)$.

Example 2.1.14. The graphs $\operatorname{CG}\left(\Omega_{p}\right)$ for all primes $p<100$ with $p \equiv 1,5$ $(\bmod 12)$ are illustrated in Figures 2.4 2.9, where the number $x$ at each vertex stands for $C(x)$, namely the cross-ratio class with respect to $x$.

Remark. For every prime $p \equiv 1(\bmod 4)$ less than $10^{5}$ and $p \neq 41, \operatorname{CG}\left(\Omega_{p}\right)$ has no bridge except its pendant edges. $\mathrm{CG}\left(\Omega_{41}\right)$ has a bridge $\{C(\chi), C(\mu)\}$. However, $\operatorname{CG}\left(\Omega_{41}\right)$ has 1-factors.


Figure 2.6: $\mathrm{CG}\left(\Omega_{53}\right)$ and $\operatorname{CG}\left(\Omega_{61}\right)$


Figure 2.7: $\mathrm{CG}\left(\Omega_{73}\right)$


Figure 2.8: $\mathrm{CG}\left(\Omega_{89}\right)$


Figure 2.9: $\mathrm{CG}\left(\Omega_{97}\right)$

### 2.1.3 Further remarks on CG graphs

In the original work of Köhler [65], he defined

$$
\begin{equation*}
K(x)=\left\{x,-1-x,-\frac{1}{1+x},-\frac{x}{1+x},-1-\frac{1}{x}, \frac{1}{x}\right\} \tag{2.9}
\end{equation*}
$$

as a vertex of "Köhler orbit graphs". By using the notation of $K(x)$, we could denote three mappings over $\mathbb{F}_{p} \cup\{\infty\}$ (cf. 2.1 ):

$$
\begin{equation*}
\tau_{\boldsymbol{A}}: x \mapsto-1-x, \tau_{\boldsymbol{B}}: x \mapsto \frac{1}{x}, \text { and } \tau_{\boldsymbol{C}}: x \mapsto \frac{-x-1}{2 x+1} . \tag{2.10}
\end{equation*}
$$

Then, one can observe that $K(x)=\left\{x^{\tau} \mid \tau \in\left\langle\tau_{\boldsymbol{A}}, \tau_{\boldsymbol{B}}\right\rangle\right\}$. In particular, $K(0)=$ $\{0,-1, \infty\}$, and $K(1)=\left\{1,-2,-2^{-1}\right\}$. It is easily seen that

$$
K(x)=-C(-x)=-\left\{-x, 1+x, \frac{1}{1+x}, \frac{x}{1+x}, 1+\frac{1}{x},-\frac{1}{x}\right\}
$$

By using $x \mapsto-x(\bmod p)(\infty \mapsto \infty)$ as the isomorphism, we can see that CG graphs and "Köhler orbit graphs" are isomorphic. Accordingly, the results obtained from Siemon's number theoretic conjecture (see [106, 108]) are also applicable to CG graphs.

Actually, it is beneficial to use $C(x)$ (cross-ratio classes) instead of $K(x)$ in our study. Especially, the complicated recursive Construction 2.3.10 for $p \equiv 1$ (mod 12) could be more reader-friendly. This is why we "redefine" those graphs to make the present work self-contained.
Remark. There is a slight difference between those two kinds of graphs. For $p \equiv 5(\bmod 12)$, Siemon [105] forbade the multiple edges, and classified the "Köhler orbit graphs" into four classes according to the congruence criteria of $p$. In contrast, LG and CG graphs allow the multiple edges, so the cases are fewer (cf. Proposition 2.1.10). Actually, the multiple edges have no effect on the constructions.

It should be pointed out that the adjacency of "Köhler orbit graphs" relies on the notion of "difference cycles", which are not purely algebraic. Even though the two families of graphs are isomorphic, the definition of $\sigma_{C}$ used for adjacencies of our CG graphs (which does not have an analogue in "Köhler orbit graphs") plays an important role in our construction below, particularly in the case when $p \equiv 1(\bmod 12)$.

So far, no efficient way is known for proving the existence of 1-factor of CG $\left(\Omega_{p}\right)$ theoretically. The proof may need to use deep mathematics in group theory and/or analytic number theory. Hence, we use computers to verify that the auxiliary graph $\mathrm{CG}^{\dagger}\left(\Omega_{p}\right)$ is 2-connected for all prime $p<10^{5}$ with $p \equiv 1$ $(\bmod 4)$.

Theorem 2.1.15. There exists an $\operatorname{AsSQS}(2 p)$ for all prime $p<10^{5}$ with $p \equiv 1$ $(\bmod 4)$.

### 2.2 Affine-invariant strictly cyclic Steiner quadruple systems over $\mathbb{Z}_{2 p}$

In this section, we begin to consider an affine-invariant sSQS (strictly cyclic Steiner quadruple system) of order $v$, provided that $v \equiv 2,10(\bmod 24)$.

Before proceeding further, we introduce some frequently used notation in this section. Let $x \in \mathscr{P}\left(\mathbb{F}_{p}\right)$. Firstly, we recall 2.5 that

$$
C(x)=\left\{x, \frac{1}{x}, \frac{x-1}{x}, \frac{x}{x-1}, \frac{1}{1-x}, 1-x\right\}
$$

is the cross-ratio class with respect to $x$. Secondly, let $\operatorname{orb}_{\boldsymbol{A C}}(x)$ denote the orbit of $x$ under the action of the subgroups $\left\langle\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{C}}\right\rangle$, i.e.,

$$
\begin{equation*}
\operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}(x)=\left\{x^{\sigma} \mid \sigma \in\left\langle\sigma_{\boldsymbol{A}}, \sigma_{\boldsymbol{C}}\right\rangle\right\}=\left\{x, 1-x, \frac{x}{2 x-1}, \frac{x-1}{2 x-1}\right\} \tag{2.11}
\end{equation*}
$$

Lastly, we simply denote

$$
\bar{x}=x^{\sigma_{C}}=\frac{1-x}{1-2 x} \quad \text { and } \quad \bar{C}(x)=C(x) \cup C(\bar{x})
$$

### 2.2.1 Block presentations

Let $v$ be a positive integer with $v \equiv 2,10(\bmod 24)$, which is the necessary condition for the existence of an $\operatorname{sSQS}(v)$ (see Proposition 1.3.1. Let $n=\frac{v}{2}$, then $n \equiv 1,5(\bmod 12)$ is odd. Since

$$
\mathbb{Z}_{2 n} \cong \mathbb{Z}_{n} \times \mathbb{Z}_{2}=\left\{(x, y) \mid x \in \mathbb{Z}_{n}, y \in \mathbb{Z}_{2}\right\}
$$

we can identify the point set $\mathbb{Z}_{v}=\mathbb{Z}_{2 n}$ with $\mathbb{Z}_{n} \times \mathbb{Z}_{2}$, and denote the point $(x, y)$ by $x_{y}$ for convenience. Addition and multiplication over $\mathbb{Z}_{n} \times \mathbb{Z}_{2}$ are defined as $x_{y}+x_{y^{\prime}}^{\prime}=\left(x+x^{\prime}\right)_{\left(y+y^{\prime}\right)}$ and $x_{y} x_{y^{\prime}}^{\prime}=\left(x x^{\prime}\right)_{\left(y y^{\prime}\right)}$, where $x+x^{\prime}, x x^{\prime}$ are reduced modulo $n$, and $y+y^{\prime}, y y^{\prime}$ are reduced modulo 2 .

For an sSQS $\left(\mathbb{Z}_{n} \times \mathbb{Z}_{2}, \mathcal{B}\right)$, let $B_{1}=\left\{a_{0}, b_{0}, c_{1}, d_{1}\right\}, B_{2}=\left\{a_{0}, b_{0}, c_{0}, d_{1}\right\}$, and $B_{3}=\left\{a_{0}, b_{0}, c_{0}, d_{0}\right\}$ be blocks in $\mathcal{B}$. By the cyclic property, $B_{1}+0_{1}=$ $\left\{a_{1}, b_{1}, c_{0}, d_{0}\right\}, B_{2}+0_{1}=\left\{a_{1}, b_{1}, c_{1}, d_{0}\right\}$, and $B_{3}+0_{1}=\left\{a_{1}, b_{1}, c_{1}, d_{1}\right\}$ are also contained in $\mathcal{B}$. Accordingly, we classify all quadruples into three types distinguished by using semicolons to separate the points.

Type I: All quadruples of the form $\left\{a_{0}, b_{0}, c_{1}, d_{1}\right\}$, simply denoted by $\{a, b ; c, d\}$, where $a \neq b$ and $c \neq d$.

Type II: All quadruples of the form $\left\{a_{0}, b_{0}, c_{0}, d_{1}\right\}$ or $\left\{a_{1}, b_{1}, c_{1}, d_{0}\right\}$, simply denoted by $\{a, b, c ; d\}$, where $a, b$, and $c$ are pairwise distinct.

Type III: All quadruples of the form $\left\{a_{0}, b_{0}, c_{0}, d_{0}\right\}$ or $\left\{a_{1}, b_{1}, c_{1}, d_{1}\right\}$, simply denoted by $\{a, b, c, d\}$, where $a, b, c$, and $d$ are pairwise distinct.

Similarly, all triples of the form $\left\{a_{0}, b_{0}, c_{0}\right\}$ or $\left\{a_{1}, b_{1}, c_{1}\right\}$ are called pure triples, simply denoted by $\{a, b, c\}$, and all triples of the form $\left\{a_{0}, b_{0}, c_{1}\right\}$ or $\left\{a_{1}, b_{1}, c_{0}\right\}$ are called mixed triples, and simply denoted by $\{a, b ; c\}$. Clearly, pure triples are contained in Type II and (or) Type III quadruples, and mixed triples are contained in Type I and (or) Type II quadruples.

Notice that

$$
\mathbb{Z}_{2 n}^{\times} \cong \mathbb{Z}_{n}^{\times} \times \mathbb{Z}_{2}^{\times}=\left\{x_{1} \mid x \in \mathbb{Z}_{n}^{\times}\right\} .
$$

For an $\operatorname{AsSQS}(2 n)$, every element $x_{1} \in \mathbb{Z}_{n}^{\times} \times \mathbb{Z}_{2}^{\times}$should be a multiplier. Hence we can simply omit the subscripts of multipliers as well. In what follows, we use these simple notation of quadruples and triples without subscripts.

### 2.2.2 Direct construction A

Let $p \equiv 1,5(\bmod 12)$ be prime, and suppose $\left(\mathbb{Z}_{2 p}, \mathcal{B}\right)$ is an AsSQS. Then, $\mathbb{Z}_{p}^{*}=\mathbb{Z}_{p}^{\times}$. For pairwise distinct elements $a, b, c, d \in \mathbb{Z}_{p}^{*}$, let $B_{1}=\{a, b ; c, d\} \in \mathcal{B}$. By the affine-invariant property,

$$
B_{1}^{\prime}:=\left(a^{-1} b-1\right)^{-1}\left(a^{-1} B_{1}-1\right)=\left\{0,1 ;(c-a)(b-a)^{-1},(d-a)(b-a)^{-1}\right\}
$$

must be contained in $\mathcal{O}_{A}\left(B_{1}\right)$. So $B_{1}^{\prime}$ can be chosen as a base block of the affine orbit. Similarly, in each affine orbit, regardless of Type I, II, or III, there must exist a block containing $\{0,1\}$. Therefore, it suffices to find base blocks of the form $\{0,1 ; a, b\},\{0,1, c ; d\}$, and $\{0,1, e, f\}$ such that all the pure triples in $\left\{\{0,1, x\} \mid x \in \mathbb{Z}_{p} \backslash\{0,1\}\right\}$ and all the mixed triples in $\left\{\{0,1 ; y\} \mid y \in \mathbb{Z}_{p}\right\}$ are covered exactly once in their affine orbits.

In order to cover the pure triple $\{0,1,-1\}$, we have the following lemma:
Lemma 2.2.1 (Type II). Let $\left(\mathbb{Z}_{2 p}, \mathcal{B}\right)$ be an AsSQS. Then there exists a Type II quadruple $\{0,1,-1 ; 0\} \in \mathcal{B}$.

Proof. Obviously, the pure triple $\{0,1,-1\}$ cannot be covered by any Type I quadruple. Now, we show that $\{0,1,-1\}$ cannot be covered by a Type III quadruple either. Assume there exists a Type III quadruple containing $\{0,1,-1\}$, say $B=\{0,1,-1, x\} \in \mathcal{B}$. Then we have $-B=\{0,-1,1,-x\} \in \mathcal{B}$, which implies $x=-x$. Hence $x=0$, in which case $B=\{0,1,-1\}$ becomes a triple. Then, we suppose there is a Type II quadruple $\{0,1,-1 ; y\} \in \mathcal{B}$. Thus $\{0,-1,1 ;-y\} \in \mathcal{B}$ which implies $y=0$. Therefore, the only possible quadruple containing $\{0,1,-1\}$ is $\{0,1,-1 ; 0\}$.

In $\mathcal{O}_{A}(\{0,1,-1 ; 0\})$, there are another two quadruples containing $\{0,1\}$, namely $\{1,0,2 ; 1\}$ and $\left\{2^{-1}, 0,1 ; 2^{-1}\right\}$. Thus it remains to consider the mixed triple $\{0,1 ; y\}$ for every $y \in \mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}\right\}$, and the pure triple $\{0,1, x\}$ for every $x \in \mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}\right\}$.
Remark. Köhler [65] and Siemon [103] used geometric representations for the triples and quadruples, namely, triangles and quadrilaterals. Lemma 2.2.1 can also follow from Siemon $103 \S 2$ concerning right triangles and equilateral triangles (cf. Section 2.1.3).

In this subsection, we intend to cover all the remaining mixed and pure triples by Type I and III quadruples, respectively. Firstly, we can obtain all Type I base blocks as follows.

Lemma 2.2.2 (Type I). Let $\left\{b_{1}, b_{2}, \ldots, b_{\frac{p-5}{4}}\right\}$ be a system of representatives of

$$
\begin{equation*}
\left\{\operatorname{orb}_{\boldsymbol{A C}}(b) \mid b \in \mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi, 1-\chi\right\}\right\} . \tag{2.12}
\end{equation*}
$$

Let $B_{b}^{(1)}=\{0,1 ; b, 1-b\}$ and

$$
\mathcal{B}_{1}=\left\{B_{b_{i}}^{(1)} \left\lvert\, i \in\left[\frac{p-5}{4}\right]\right.\right\} \cup\left\{B_{\chi}^{(1)}\right\}
$$

Then, $\bigcup_{B \in \mathcal{B}_{1}} \mathcal{O}_{A}(B)$ covers the mixed triple $\{0,1 ; y\}$ exactly once for every $y \in \mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}\right\}$.

Proof. There are only two quadruples containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{b}^{(1)}\right)$, namely $B_{b}^{(1)}=\{0,1 ; b, 1-b\}$ and $\left\{\frac{-b}{1-2 b}, \frac{1-b}{1-2 b} ; 0,1\right\}$. Hence, all the mixed triples contained in $\mathcal{O}_{A}\left(B_{b}^{(1)}\right)$ are $\left\{\{0,1 ; y\} \mid y \in \operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}(b)\right\}$. In particular, $\operatorname{orb}_{\boldsymbol{A C}}(0)=$ $\{0,1\}, \operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(2^{-1}\right)=\left\{2^{-1}, \infty\right\}$, and $\operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}(\chi)=\{\chi, 1-\chi\}$. Thus,

$$
\biguplus_{i=1}^{\frac{p-5}{4}} \operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i}\right) \cup \operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}(\chi)=\bigcup_{b \in \mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}\right\}} \operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}(b)=\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}\right\}
$$

which implies that each mixed triple containing $\{0,1\}$ appears exactly once in $\bigcup_{B \in \mathcal{B}_{1}} \mathcal{O}_{A}(B)$.

Remark. Köhler 65] separated the "Köhler graph" into two components, $\mathrm{KG}_{1}$ and $\mathrm{KG}_{2}$, where $\mathrm{KG}_{1}$ was proved to have a 1 -factor by Siemon [103]. Lemma 2.2.2 can also follow from Siemon 103 $\S 3$ applied to $\mathrm{KG}_{1}$ (cf. Section 2.1.3).

It has been shown that $\Omega_{p}=\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}\right\}$ can be partitioned into cross-ratio classes. Moreover, each edge in the graph $\mathrm{CG}\left(\Omega_{p}\right)$ can be written as $\{C(a), C(\bar{a})\}$ for some $a \in \Omega_{p}$, where $C(a)$ is the cross-ratio class with respect to $a$. Suppose $\operatorname{CG}\left(\Omega_{p}\right)$ has a 1-factor, then the edge set of the 1 -factor gives a partition of $\Omega_{p}$ into pairs of cross-ratio classes. Moreover, when $p \equiv 1(\bmod 12)$, by Proposition 2.1 .10 (ii), there is a pendant edge incident with $C(\xi)$, which must be contained in the 1 -factor. We propose the base block corresponding to this pendant edge as follows:

Lemma 2.2.3 (Type $\left.^{\mathrm{III}}{ }^{\xi}\right)$. Suppose $p \equiv 1(\bmod 12)$. If $\mathrm{CG}\left(\Omega_{p}\right)$ has a 1 -factor $F$ containing the edge $\{C(\xi), C(\bar{\xi})\}$, where $\xi$ is a root of $\xi^{2}-\xi+1=0$ over $\mathbb{Z}_{p}$. Let $B_{\xi}^{(3)}=\{0,1, \xi, \bar{\xi}\}$. Then, $\mathcal{O}_{A}\left(B_{\xi}^{(3)}\right)$ covers the pure triple $\{0,1, x\}$ exactly once for every $x \in \bar{C}(\xi)$.

Proof. All the quadruples containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{\xi}^{(3)}\right)$ can be determined explicitly, namely, $\left\{0,1, \xi, \frac{1}{1+\xi^{-1}}\right\},\left\{0, \xi^{-1}, 1, \frac{1}{1+\xi}\right\},\left\{0,1+\xi^{-1}, 1+\xi, 1\right\}$, and $\left\{1,-\xi^{-1},-\xi, 0\right\}$. Note that

$$
C(\bar{\xi})=\left\{\frac{1}{1+\xi^{-1}}, 1+\xi^{-1},-\xi^{-1},-\xi, 1+\xi, \frac{1}{1+\xi}\right\} .
$$

Hence, the pure triple $\{0,1, x\}$ for every $x \in \bar{C}(\xi)=C(\xi) \cup C(\bar{\xi})$ appears exactly once in $\mathcal{O}_{A}\left(B_{\xi}^{(3)}\right)$.

The graph $\operatorname{CG}\left(\Omega_{p}\right)$ has $\frac{p-1}{6}$ vertices if $p \equiv 1(\bmod 12)$, and $\frac{p-5}{6}$ vertices if $p \equiv 5(\bmod 12)$. Hence, besides the previously discussed edge $\{C(\xi), C(\bar{\xi})\}$ when $p \equiv 1(\bmod 12)$, a 1 -factor of $\mathrm{CG}\left(\Omega_{p}\right)$ has $\ell_{p}$ edges, where

$$
\ell_{p}=\left\{\begin{array}{lll}
\frac{p-5}{12}, & \text { if } p \equiv 5 & (\bmod 12)  \tag{2.13}\\
\frac{p-13}{12}, & \text { if } p \equiv 1 \quad(\bmod 12)
\end{array}\right.
$$

We can obtain all the other Type III base blocks as follows:
Lemma 2.2.4 (Type III). Assume that $\operatorname{CG}\left(\Omega_{p}\right)$ has a 1-factor $F$. Let $a_{1}, a_{2}$, $\ldots, a_{\ell_{p}}$ be elements in $\Omega_{p}$ with $a_{i} \notin C(\xi) \cup C(\bar{\xi})$ for each $i \in\left[\ell_{p}\right]$, such that
$\left\{\left\{C\left(a_{1}\right), C\left(\overline{a_{1}}\right)\right\},\left\{C\left(a_{2}\right), C\left(\overline{a_{2}}\right)\right\}, \ldots,\left\{C\left(a_{\ell_{p}}\right), C\left(\overline{a_{\ell_{p}}}\right)\right\}\right\}=E(F) \backslash\{C(\xi), C(\bar{\xi})\}$,
where $E(F)$ is the edge set of $F$. Let $B_{a_{i}}^{(3)}=\left\{0,1, a_{i}, 1-a_{i}\right\}$ and $\mathcal{B}_{3}=$ $\left\{B_{a_{i}}^{(3)} \mid i \in\left[\ell_{p}\right]\right\}$. Then $\bigcup_{B \in \mathcal{B}_{3}} \mathcal{O}_{A}(B)$ covers the pure triple $\{0,1, x\}$ exactly once for every $x \in \Omega_{p} \backslash \bar{C}(\xi)$.

Proof. For $a \notin \bar{C}(\xi)$, all the quadruples containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{a}^{(3)}\right)$ are $\{0,1, a, 1-a\},\left\{0, \frac{1}{a}, 1, \frac{\bar{a}}{\bar{a}-1}\right\},\left\{1, \frac{a-1}{a}, 0, \frac{1}{1-\bar{a}}\right\},\left\{0, \frac{1}{1-a}, \frac{\bar{a}-1}{\bar{a}}, 1\right\},\left\{1, \frac{a}{a-1}, \frac{1}{\bar{a}}, 0\right\}$, and $\{\bar{a}, 1-\bar{a}, 1,0\}$. Therefore, the pure triple $\{0,1, x\}$ appears in $\mathcal{O}_{A}\left(B_{a}^{(3)}\right)$ for every $x \in \bar{C}(a)$. Furthermore, each pure triple $\{0,1, x\}$ appears exactly once if $|C(a)|=|C(\bar{a})|=6$ holds. In fact, as shown in 2.6), all the crossratio classes are of size 6 except $C(-1)=\left\{-1,2,2^{-1}\right\}, C(0)=\{0,1, \infty\}$, and $C(\xi)=\{\xi, 1-\xi\}$.

To summarize, we propose Theorem 2.2 .5 and Construction 2.2 .6 which directly follow from Lemmas 2.2.1, 2.2.2, 2.2.3, and 2.2.4. Theorem 2.2.5 indicates that $\operatorname{AsSQS}(v)$, which has a stronger "symmetry", does not require stronger criteria than $\operatorname{sSQS}(v)$.

Theorem 2.2.5. Let $p \equiv 1(\bmod 4)$ be a prime. If $\mathrm{CG}\left(\Omega_{p}\right)$ has a 1-factor, then there exists an AsSQS(2p).

Table 2.2: The base blocks of an $\operatorname{AsSQS}^{A}(2 p)$ for $p \equiv 5(\bmod 12)$

| Type | Base blocks | \# Base blocks | \# Cyclic orbits | Lemmas |
| ---: | :--- | :--- | :--- | :--- |
| $\mathrm{I}^{\prime}$ | $\{0,1 ; \chi, 1-\chi\}$ | 1 | $\frac{p-1}{4}$ | Lemma |
| I | $\left\{0,1 ; b_{i}, 1-b_{i}\right\}$ | $i \in\left[\frac{p-5}{4}\right]$ | $\frac{p-1}{2}$ | Lemma |
| $\mathrm{II}^{\prime}$ | $\{0,1,-1 ; 0\}$ | 1 | $\frac{p-1}{2.2 .2}$ |  |
| III | $\left\{0,1, a_{i}, 1-a_{i}\right\}$ | $i \in\left[\frac{p-5}{12}\right]$ | $\frac{p-1}{2}$ | Lemma |
| Total |  | $\frac{p+1}{3}$ | $\frac{(p-1)(2 p-1)}{12}$ | Lemma |

Table 2.3: The base blocks of an $\operatorname{AsSQS}^{A}(2 p)$ for $p \equiv 1(\bmod 12)$

| Type | Base blocks | \# Base blocks | \# Cyclic orbits | Lemmas |
| ---: | :--- | :--- | :--- | :--- |
| $\mathrm{I}^{\prime}$ | $\{0,1 ; \chi, 1-\chi\}$ | 1 | $\frac{p-1}{4}$ | Lemma 2.2 .2 |
| I | $\left\{0,1 ; b_{i}, 1-b_{i}\right\}$ | $i \in\left[\frac{p-5}{4}\right]$ | $\frac{p-1}{2}$ | Lemma |
| $\mathrm{II}^{\prime}$ | $\{0,1,-1 ; 0\}$ | 1 | $\frac{p-1}{2.2}$ |  |
| $\mathrm{III}^{\xi}$ | $\{0,1, \xi, \bar{\xi}\}$ | 1 | Lemma | $\frac{p-1}{2.2 .1}$ |
| III | $\left\{0,1, a_{i}, 1-a_{i}\right\}$ | $i \in\left[\frac{p-13}{12}\right]$ | $\frac{p-1}{2}$ | Lemma |
| Total |  | $\frac{p+2}{3}$ | $\frac{(p-1)(2 p-1)}{12}$ | Lemma |
| To.2.4 |  |  |  |  |

Construction 2.2.6. If $\mathrm{CG}\left(\Omega_{p}\right)$ has a 1 -factor, choose $a_{1}, a_{2}, \ldots, a_{\left\lfloor\frac{p}{12}\right\rfloor}$ as in Lemma 2.2.4 Choose $b_{1}, b_{2}, \ldots, b_{\frac{p-1}{4}}$ as in Lemma 2.2.2 Then all the base blocks of $\operatorname{AsSQS}(2 p)$ are given as follows:
(i) For $p \equiv 1(\bmod 12)$,

Type I, $\quad\left\{0,1 ; b_{i}, 1-b_{i}\right\}$, for $i \in\left[\frac{p-1}{4}\right]$,
Type II' $^{\prime}, \quad\{0,1,-1 ; 0\}$,
Type $\mathrm{III}^{\xi}, \quad\{0,1, \xi, \bar{\xi}\}$,
Type III, $\left\{0,1, a_{i}, 1-a_{i}\right\}$, for $i \in\left[\frac{p-13}{12}\right], a_{i} \notin \bar{C}(\xi)$,
where $\xi$ is a root of $x^{2}-x+1=0$ over $\mathbb{Z}_{p}$.
(ii) For $p \equiv 5(\bmod 12)$,

Type I, $\left\{0,1 ; b_{i}, 1-b_{i}\right\}$, for $i \in\left[\frac{p-1}{4}\right]$,
Type II $^{\prime}, \quad\{0,1,-1 ; 0\}$,
Type III, $\left\{0,1, a_{i}, 1-a_{i}\right\}$, for $i \in\left[\frac{p-5}{12}\right]$.
In Table 2.2 and Table 2.3, we summarize the number of base blocks of each type (in the column with header "\# Base blocks"), and the numbers of cyclic orbits contained in the affine orbit of a base block $B$ of a given type (in the column with header "\# Cyclic orbits"), that is, $\frac{\left|\mathcal{O}_{A}(B)\right|}{2 p}$.

Moreover, we denote an $\operatorname{AsSQS}(2 p)$ obtained from Constructions 2.2.6 by $\operatorname{AsSQS}^{A}(2 p)$.

By utilizing Theorem 2.1.12 and Lemma 2.1.13, respectively, we have the following as corollaries of Theorem 2.2.5

Corollary 2.2.7. For prime $p \equiv 1(\bmod 4)$ and $p \not \equiv 1,49(\bmod 60)$, if $\mathrm{CG}\left(\Omega_{p}\right)$ is bridgeless, then there exists an $A s S Q S^{A}(2 p)$.

Corollary 2.2.8. There exists an $\operatorname{AsSQ} S^{A}(2 p)$ for all primes $p<10^{5}$ satisfying $p \equiv 1(\bmod 4)$.

Example 2.2.9. Let $p=17$. Then $\mathcal{B}_{1}$ consists of

$$
\{0,1 ; 7,11\},\{0,1 ; 4,14\},\{0,1 ; 5,13\},\{0,1 ; 6,12\}
$$

and $\mathcal{B}_{3}$ consists of

$$
\{0,1,3,15\}
$$

where $\chi=7$ (cf. Figure 2.4).
Example 2.2.10. Let $p=29$. Then $\mathcal{B}_{1}$ consists of

$$
\begin{gathered}
\{0,1 ; 9,21\},\{0,1 ; 13,17\},\{0,1 ; 12,18\},\{0,1 ; 10,20\}, \\
\{0,1 ; 14,16\},\{0,1 ; 11,19\},\{0,1 ; 7,25\},
\end{gathered}
$$

and $\mathcal{B}_{3}$ consists of

$$
\{0,1,3,27\},\{0,1,11,19\}
$$

which correspond to the edges $\{C(3), C(9)\}$ and $\{C(4), C(5)\}$ of $\mathrm{CG}\left(\Omega_{29}\right)$, respectively, where $\chi=9$ (cf. Figure 2.5.

Example 2.2.11. Let $p=41$. Then $\mathcal{B}_{1}$ consists of

$$
\begin{gathered}
\{0,1 ; 5,37\},\{0,1 ; 14,28\},\{0,1 ; 17,25\},\{0,1 ; 18,24\}, \\
\{0,1 ; 8,34\},\{0,1 ; 10,32\},\{0,1 ; 15,27\},\{0,1 ; 20,22\}, \\
\{0,1 ; 13,29\},\{0,1 ; 19,23\},
\end{gathered}
$$

and $\mathcal{B}_{3}$ consists of

$$
\{0,1,3,39\},\{0,1,4,38\},\{0,1,10,32\}
$$

which correspond to the edges $\{C(3), C(13)\},\{C(4), C(12)\}$, and $\{C(5), C(6)\}$ of $\mathrm{CG}\left(\Omega_{41}\right)$, respectively, where $\chi=5$ (cf. Figure 2.5).

### 2.2.3 Direct construction B

Yoshikawa [129] proposed an idea for constructing AsSQS(2p) by using Type II quadruples as much as possible, and showed that such $\operatorname{AsSQS}(2 p)$ exists for all primes $17 \leq p<100$ satisfying $p \equiv 1(\bmod 4)$ by computer search. In this subsection, we present Yoshikawa's idea somewhat differently and provide a combinatorial criterion for those constructions.

Let $p \equiv 1,5(\bmod 12)$ be prime and suppose $\left(\mathbb{Z}_{2 p}, \mathcal{B}\right)$ is an AsSQS without Type III quadruples. If we color the point $x \in \mathbb{Z}_{2 p}$ in red if $x \equiv 0(\bmod 2)$ and in blue if $x \equiv 1(\bmod 2)$, it is clear that Type III quadruples are monochromatic and the other two types are not. By assigning two colors to all the points, if an SQS does not have any monochromatic quadruple, then the SQS is said to be 2-chromatic (see [59, 92]). Hence, an AsSQS having no Type III quadruples must be 2-chromatic. Note that, under the assumptions in this subsection, Lemma 2.2.1 remains true, that is, $\{0,1,-1 ; 0\} \in \mathcal{B}$.

Lemma 2.2.12 (Type I). There exists exactly one affine orbit of Type I quadruples, say $\mathcal{O}_{A}(B)$, where $B=\{0,1 ; \chi, 1-\chi\}$, and $\chi=\frac{1+\sqrt{-1}}{2}$ is a root of $2 \chi^{2}-2 \chi+1=0$ over $\mathbb{Z}_{p}$.

Proof. The total numbers of pure triples $\{0,1, x\}$ and mixed triples $\{0,1 ; y\}$ are $p-2$ and $p$, respectively. Under the assumption that there is no Type III quadruple, there must be exactly two mixed triples not covered by Type II quadruples, say $\left\{0,1 ; y_{1}\right\}$ and $\left\{0,1 ; y_{2}\right\}$. Let $B=\left\{0,1 ; y_{1}, y_{2}\right\}$ be a quadruple in $\mathcal{B}$. Then there should not exist any other triple containing $\{0,1\}$ in $\mathcal{O}_{A}(B)$. Note that $1-B=\left\{1,0 ; 1-y_{1}, 1-y_{2}\right\} \in \mathcal{O}_{A}(B)$. Then we have $\left\{y_{1}, y_{2}\right\}=\{1-$ $\left.y_{1}, 1-y_{2}\right\}$ which implies $y_{1}+y_{2}=1$. For $B=\left\{0,1 ; y_{1}, 1-y_{1}\right\}$, similarly to the proof of Lemma 2.2.2, $\operatorname{orb}_{\boldsymbol{A C}}\left(y_{1}\right)$ should be of size two (cf. 2.11). Thus we have $y_{1} \in\{0,1\}$ or $y_{1} \in\{\chi, 1-\chi\}$, in which the former contradicts Lemma 2.2.1.

Let $B=\{0,1, x ; y\}$ be a Type II quadruple. Then there are six quadruples containing $\{0,1\}$ in $\mathcal{O}_{A}(B)$, namely,

$$
\begin{gathered}
\{0,1, x ; y\},\{1,0,1-x ; 1-y\},\left\{\frac{1}{1-x}, 0,1 ; \frac{1-y}{1-x}\right\} \\
\left\{\frac{x}{x-1}, 1,0 ; \frac{x-y}{x-1}\right\},\left\{1, \frac{x-1}{x}, 0 ; \frac{x-y}{x}\right\}, \text { and }\left\{0, \frac{1}{x}, 1 ; \frac{y}{x}\right\} .
\end{gathered}
$$

Clearly, $\mathcal{O}_{A}(B)$ covers the pure triple $\{0,1, a\}$ for every $a \in C(x)$ and the mixed triple $\{0,1 ; b\}$ for every $b \in H(x, y)$, where $C(x)$ is the cross-ratio class with respect to $x$ and

$$
\begin{equation*}
H(x, y)=\left\{y, \frac{y}{x}, \frac{x-y}{x}, \frac{x-y}{x-1}, \frac{1-y}{1-x}, 1-y\right\} \tag{2.14}
\end{equation*}
$$

Lemma 2.2.13 $\left(\right.$ Type $\left.\mathrm{II}^{\xi}\right)$. For $p \equiv 1(\bmod 12),\{0,1, \xi ; \bar{\xi}\} \in \mathcal{B}$ holds, where $\xi$ is a root of $\xi^{2}-\xi+1=0$ over $\mathbb{Z}_{p}$, and $\bar{\xi}=\frac{\xi}{\xi+1}$.

Proof. By 2.6 , $C(\xi)=\{\xi, 1-\xi\}$ is of size two. Under the assumption that there is no Type III quadruple, we suppose $B=\{0,1, \xi ; y\} \in \mathcal{B}$. All the quadruples containing $\{0,1, \xi\}$ in $\mathcal{O}_{A}(B)$ should be identical, that is, $\{0,1, \xi ; y\}$, $\{0,1, \xi ; \xi(1-y)\}$, and $\left\{0,1, \xi ; 1-\xi^{-1} y\right\}$ are identical. Hence, we have $y=$ $\frac{\xi}{\xi+1}$.

Let

$$
\begin{align*}
& \Omega_{p}^{*}= \begin{cases}\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}, \xi, 1-\xi\right\}, & \text { for } p \equiv 1(\bmod 12) \\
\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}\right\}, & \text { for } p \equiv 5(\bmod 12)\end{cases}  \tag{2.15}\\
& \Lambda_{p}= \begin{cases}\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi, 1-\chi, \bar{\xi}, 1-\bar{\xi}\right\}, & \text { for } p \equiv 1(\bmod 12) \\
\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi, 1-\chi\right\}\end{cases} \tag{2.16}
\end{align*}
$$

It is clear that $\left|\Omega_{p}^{*}\right|=\left|\Lambda_{p}\right|$ is divisible by 6 . By combining Lemmas 2.2.1 2.2.12 and 2.2.13, it is straightforward to get the following:

Lemma 2.2.14 (Type II). An AsSQS(2p) having no Type III quadruples exists if and only if there exists $\left(x_{i}, y_{i}\right) \in \Omega_{p}^{*} \times \Lambda_{p}$ for each $i \in\left[\frac{\left|\Lambda_{p}\right|}{6}\right]$, such that

$$
\biguplus_{i=1}^{\left|\Lambda_{p}\right| / 6} C\left(x_{i}\right)=\Omega_{p}^{*} \quad \text { and } \quad \bigcup_{i=1}^{\left|\Lambda_{p}\right| / 6} H\left(x_{i}, y_{i}\right)=\Lambda_{p}
$$

It is easily seen from 2.6 and 2.15 that $\Omega_{p}^{*}$ can be partitioned into crossratio classes of size 6 . Next, we intend to partition $\Lambda_{p}$ into $H(x, y)$. Let $z=\frac{y}{x}$, then $H(x, y)$ can be rewritten as

$$
\begin{equation*}
H\left(\frac{y}{z}, y\right)=\left\{y, z, 1-z, \frac{y(1-z)}{y-z}, \frac{z(1-y)}{z-y}, 1-y\right\} \tag{2.17}
\end{equation*}
$$

For $y, z \in \mathbb{Z}_{p}$, we define a mapping $\diamond: \mathbb{Z}_{p} \times \mathbb{Z}_{p} \rightarrow \mathbb{Z}_{p} \cup\{\infty\}$ as a binary operator, such that $y \diamond z=\frac{y(1-z)}{y-z}$, and $y \diamond z=\infty$ if $y=z$. The following properties can be easily derived:

Proposition 2.2 .15 . For any distinct $y, z \in \mathbb{Z}_{p} \backslash\{0,1\}$, the following hold:
(i) $(1-y) \diamond(1-z)=z \diamond y$;
(ii) $y \diamond z+z \diamond y=1$;
(iii) The solutions of $y \diamond x=z$ and $x \diamond y=z$ with respect to $x$ are $x=y \diamond z$ and $x=y \diamond(1-z)$, respectively.

In order to establish a hypergraph representation and convert the partition problem to a similar hypergraph factor problem (cf. the graph $\operatorname{CG}\left(\Omega_{p}\right)$ ), we need some notation and propositions. For a given $x \in \mathbb{Z}_{p}$, denote

$$
\llbracket x \rrbracket= \begin{cases}x \quad(\bmod p) & \text { if } x \in\left\{1,2, \ldots, \frac{p+1}{2}\right\},  \tag{2.18}\\ 1-x \quad(\bmod p) & \text { if } x \in\left\{\frac{p+3}{2}, \frac{p+5}{2}, \ldots, p-1,0\right\} .\end{cases}
$$

For a subset $X \subseteq \mathbb{Z}_{p}$, denote

$$
\llbracket X \rrbracket=\{\llbracket x \rrbracket \mid x \in X\} .
$$

For distinct $y, z \in \mathbb{Z}_{p}$, define

$$
\begin{equation*}
J(y, z)=\llbracket H\left(\frac{y}{z}, y\right) \rrbracket=\{\llbracket y \rrbracket, \llbracket z \rrbracket, \llbracket y \diamond z \rrbracket\} \tag{2.19}
\end{equation*}
$$

It is easy to observe that $\biguplus_{i=1}^{\left|\Lambda_{p}\right| / 6} H\left(x_{i}, y_{i}\right)=\Lambda_{p}$ if and only if

$$
\left|\Lambda_{p}\right| / 6
$$

$$
\biguplus_{i=1} J\left(y_{i}, z_{i}\right)=\llbracket \Lambda_{p} \rrbracket,
$$

where $z_{i}=y_{i} / x_{i}$.
Lemma 2.2.16. Let

$$
\begin{equation*}
\mathcal{B}=\left\{J(y, z) \mid y, z \in \mathbb{Z}_{p} \backslash\{0,1\}, y \neq z\right\} \tag{2.20}
\end{equation*}
$$

(i) For any distinct $y, z \in \mathbb{Z}_{p} \backslash\{0,1\}$ with $y+z \neq 1$, the pair $\{\llbracket y \rrbracket, \llbracket z \rrbracket\}$ appears exactly $\lambda$ times in $\mathcal{B}$, where $\lambda= \begin{cases}1 & \text { if } 2^{-1} \in\{y, z\}, \\ 2 & \text { otherwise. }\end{cases}$
(ii) For any $B \in \mathcal{B}, \#\left\{\left.C\left(\frac{u}{v}\right) \right\rvert\, u, v \in \mathbb{Z}_{p} \backslash\{0,1\}, J(u, v)=B\right\} \leq 2$.

Proof. (i) The pair $\{\llbracket y \rrbracket, \llbracket z \rrbracket\}$ is contained in $J(y, z)$ and $J(y, 1-z)$, whose third elements are determined by $w_{1}=\llbracket y \diamond z \rrbracket$ and $w_{2}=\llbracket y \diamond(1-z) \rrbracket$, respectively. Suppose $w$ satisfies $J(a, w)=\{\llbracket a \rrbracket, \llbracket w \rrbracket, \llbracket b \rrbracket\}$ or $J(w, a)=$ $\{\llbracket w \rrbracket, \llbracket a \rrbracket, \llbracket b \rrbracket\}$, where $\{a, b\} \in\{y, 1-y\} \times\{z, 1-z\} \cup\{z, 1-z\} \times\{y, 1-y\}$. By Proposition 2.2.15 (iii), the set of all possible values of $w$ is given by

$$
W=\{a \diamond b \mid\{a, b\} \in\{y, 1-y\} \times\{z, 1-z\} \cup\{z, 1-z\} \times\{y, 1-y\}\}
$$

By using Proposition 2.2 .15 (iii) and (i) successively, we have

$$
\begin{aligned}
\llbracket W \rrbracket & =\{\llbracket a \diamond b \rrbracket \mid\{a, b\} \in\{y, 1-y\} \times\{z, 1-z\}\} \\
& =\{\llbracket y \diamond z \rrbracket, \llbracket y \diamond(1-z) \rrbracket\}=\left\{w_{1}, w_{2}\right\} .
\end{aligned}
$$

Moreover, $w_{1}=w_{2}$ if and only if $z=2^{-1}$ or $y=2^{-1}$.
(ii) For a given $B=\{\llbracket y \rrbracket, \llbracket z \rrbracket, \llbracket y \diamond z \rrbracket\}$, all the possible unordered pairs $(u, v)$ chosen from $\{y, 1-y, z, 1-z, y \diamond z, z \diamond y\}$ satisfying $J(u, v)=J(y, z)$ are $(y, z),(1-y, 1-z),(y, y \diamond z),(1-y, z \diamond y),(z, z \diamond y)$, and $(1-z, y \diamond z)$. By calculating the cross-ratio classes with respect to the quotients of all these pairs, we have $C\left(\frac{y}{z}\right)=C\left(\frac{1-y}{z \diamond y}\right)=C\left(\frac{1-z}{y \diamond z}\right)$ and $C\left(\frac{1-y}{1-z}\right)=C\left(\frac{y}{y \diamond z}\right)=$ $C\left(\frac{z}{z \diamond y}\right)$, which completes the proof.

Remark. From the viewpoint of hypergraphs, we consider the vertex-induced sub-hypergraph of $\left(\llbracket \mathbb{Z}_{p} \rrbracket \backslash\{0,1\}, \mathcal{B}\right)$ on $\llbracket \Lambda_{p} \rrbracket$, whose edge set (treated as a multiset) is

$$
\mathcal{B}^{*}=\left\{B^{*}=B \backslash \llbracket \mathbb{Z}_{p} \backslash \Lambda_{p} \rrbracket\left|B \in \mathcal{B},\left|B^{*}\right| \geq 2\right\}\right.
$$

Since $2^{-1} \notin \Lambda_{p}$, by Lemma 2.2.16. $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}^{*}\right)$ is a $\left(\frac{\left|\Lambda_{p}\right|}{2},\{2,3\}, 2\right)$ pairwise balanced design (PBD).

Definition 2.2.17. Let $\mathcal{B}^{\triangle}$ denote the collection of all triples in $\mathcal{B}^{*}$, i.e.,

$$
\begin{equation*}
\mathcal{B}^{\triangle}=\left\{B^{*}=B \backslash \llbracket \mathbb{Z}_{p} \backslash \Lambda_{p} \rrbracket\left|B \in \mathcal{B},\left|B^{*}\right|=3\right\}\right. \tag{2.21}
\end{equation*}
$$

Let

$$
\begin{equation*}
\mathcal{C}=\left\{C(x) \mid x \in \Omega_{p}^{*}\right\} \tag{2.22}
\end{equation*}
$$

Let $\gamma_{0}$ denote a mapping from $\mathcal{B}^{\triangle}$ to $2^{\mathcal{C}}$ such that

$$
\gamma_{0}(B)=\left\{\left.C\left(\frac{u}{v}\right) \right\rvert\, u, v \in \Lambda_{p}, J(u, v)=B\right\} .
$$

Then, $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}^{\triangle}, \gamma_{0}\right)$ is an edge-colored hypergraph.
It follows from Lemma 2.2 .16 (iii) that $\left|\gamma_{0}(B)\right| \leq 2$ for every $B \in \mathcal{B}^{\triangle}$. Thus, it suffices to investigate an equivalent edge-colored hypergraph whose color-set is $\mathcal{C}$.
Definition 2.2.18. Let $\mathcal{B}_{2}^{\triangle}$ be the multiset obtained by doubling all triples in $\mathcal{B}^{\triangle}$, i.e., $\mathcal{B}_{2}^{\triangle}=\mathcal{B}^{\triangle} \uplus \mathcal{B}^{\triangle}$. Let $\gamma$ be a mapping from $\mathcal{B}_{2}^{\triangle}$ to $\mathcal{C}$, such that for given $B=B^{\prime}=\{\llbracket y \rrbracket, \llbracket z \rrbracket, \llbracket y \diamond z \rrbracket\} \in \mathcal{B}_{2}^{\triangle}$,

$$
\gamma(B)=C\left(\frac{y}{z}\right) \quad \text { and } \quad \gamma\left(B^{\prime}\right)=C\left(\frac{1-y}{1-z}\right)
$$

where $\left\{C\left(\frac{y}{z}\right), C\left(\frac{1-y}{1-z}\right)\right\}=\gamma_{0}(B)$ for $B \in \mathcal{B}^{\triangle}$. Then $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}_{2}^{\triangle}, \gamma\right)$ is an edgecolored hypergraph.

Suppose $\mathcal{F} \subseteq \mathcal{B}^{\triangle}$ is a 1-factor of $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}_{2}^{\triangle}, \gamma\right)$. If $\gamma\left(F_{1}\right) \neq \gamma\left(F_{2}\right)$ for any distinct $F_{1}, F_{2} \in \mathcal{F}$, then $\mathcal{F}$ is said to be a rainbow 1-factor. To sum it all up, we propose Theorem 2.2.19 and Construction 2.2.20, using Lemmas 2.2.1, 2.2.12. 2.2.13 and 2.2.14

Theorem 2.2.19. An $A s S Q S(2 p)$ having no Type III quadruples exists if the edge-colored hypergraph $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}_{2}^{\triangle}, \gamma\right)$ has a rainbow 1-factor.

Proof. A 1-factor of $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}_{2}^{\triangle}, \gamma\right)$ gives rise to a partition of $\llbracket \Lambda_{p} \rrbracket$ into triples of the form $J(y, z)$. Moreover, it is clear that $\mathcal{C}$ is a partition of $\Omega_{p}^{*}$, which is also the color-set of $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}_{2}^{\triangle}, \gamma\right)$. Since $\frac{\left|\Omega_{p}^{*}\right|}{6}=\frac{\left|\llbracket \Lambda_{p} \rrbracket\right|}{3}$, every edge in a rainbow 1 -factor has a distinct color in $\mathcal{C}$. Hence, a rainbow 1 -factor forms a partition of $\Omega_{p}^{*}$. By Lemma 2.2.14, an $\operatorname{AsSQS}(2 p)$ having no Type III quadruples can be constructed as claimed.
Construction 2.2.20. When $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}_{2}^{\triangle}, \gamma\right)$ has a rainbow 1-factor with edge set

$$
\mathcal{F}=\left\{B_{i}=J\left(y_{i}, z_{i}\right) \left\lvert\, i \in\left[\frac{\left|\Lambda_{p}\right|}{6}\right]\right.\right\}
$$

such that

$$
\bigoplus_{i=1}^{\left|\Lambda_{p}\right| / 6} C\left(\frac{y_{i}}{z_{i}}\right)=\Omega_{p}^{*}
$$

Table 2.4: The base blocks of an $\operatorname{AsSQS}^{B}(2 p)$ for $p \equiv 5(\bmod 12)$

| Type | Base block | \# Base blocks | \# Cyclic orbits | Lemmas |
| ---: | :--- | :--- | :--- | :--- |
| I' $^{\prime}$ | $\{0,1 ; \chi, 1-\chi\}$ | 1 | $\frac{p-1}{4}$ | Lemma 2.2.12 |
| II $^{\prime}$ | $\{0,1,-1 ; 0\}$ | 1 | $\frac{p-1}{2}$ | Lemma 2.2 .1 |
| II | $\left\{0,1, a_{i} ; b_{i}\right\}$ | $i \in\left[\frac{p-5}{6}\right]$ | $p-1$ | Lemma 2.2 .14 |
| Total |  | $\frac{p+7}{6}$ | $\frac{(p-1)(2 p-1)}{12}$ |  |

Table 2.5: The base blocks of an $\operatorname{AsSQS}^{B}(2 p)$ for $p \equiv 1(\bmod 12)$

| Type | Base block | \# Base blocks | \# Cyclic orbits | Lemmas |
| ---: | :--- | :--- | :--- | :--- |
| $\mathrm{I}^{\prime}$ | $\{0,1 ; \chi, 1-\chi\}$ | 1 | $\frac{p-1}{4}$ | Lemma 2.2.12 |
| $\mathrm{II}^{\prime}$ | $\{0,1,-1 ; 0\}$ | 1 | $\frac{p-1}{2}$ | Lemma |
| $\mathrm{II}^{\xi}$ | $\{0,1, \xi ; \bar{\xi}\}$ | 1 | $\frac{p-1}{3}$ | Lemma |
| II | $\left\{0,1, a_{i} ; b_{i}\right\}$ | $i \in\left[\frac{p-7}{6}\right]$ | $p-1$ | Lemma |
| Total |  | $\frac{p+11}{6}$ | $\frac{(p-1)(2 p-1)}{12}$ |  |

let $a_{i}=\frac{y_{i}}{z_{i}}$ and $b_{i}=y_{i}$ for $i \in\left[\frac{\left|\Lambda_{p}\right|}{6}\right]$. The base blocks of an $\operatorname{AsSQS}(2 p)$ are given as follows:
(i) For $p \equiv 1(\bmod 12)$,

Type I, $\quad\{0,1 ; \chi, 1-\chi\}$,
Type $\mathrm{II}^{\prime}, \quad\{0,1,-1 ; 0\}$,
Type $\mathrm{II}^{\xi}, \quad\{0,1, \xi ; \bar{\xi}\}$,
Type II, $\left\{0,1, a_{i} ; b_{i}\right\}$ for $i \in\left[\frac{p-7}{6}\right]$,
where $\xi$ is a root of $\xi^{2}-\xi+1=0$ over $\mathbb{Z}_{p}$.
(ii) For $p \equiv 5(\bmod 12)$,

Type I, $\quad\{0,1 ; \chi, 1-\chi\}$,
Type $\mathrm{II}^{\prime}, \quad\{0,1,-1 ; 0\}$,
Type II, $\quad\left\{0,1, a_{i} ; b_{i}\right\}$ for $i \in\left[\frac{p-5}{6}\right]$,
where $\chi$ is a root of $2 \chi^{2}-2 \chi+1=0$ over $\mathbb{Z}_{p}$.
In Table 2.4 and Table 2.5, we summarize the number of base blocks of each type (in the column with the header "\# Base blocks"), and the numbers of cyclic orbits covered by the affine orbit of a given base block $B$ in each type (in the column with the header "\# Cyclic orbits"), that is, $\frac{\left|\mathcal{O}_{A}(B)\right|}{2 p}$.

We denote an $\operatorname{AsSQS}(2 p)$ obtained from Construction 2.2 .20 by $\operatorname{AsSQS}^{B}(2 p)$.

Remark. As shown in Tables 2.2, 2.3, 2.4, and 2.5, the number of base blocks of an $\operatorname{AsSQS}^{B}(2 p)$ is approximately half of $\operatorname{AsSQS}^{A}(2 p)$ for a fixed $p$. However, the total numbers of cyclic orbits are the same. It is clear that an $\operatorname{AsSQS}^{A}(2 p)$ and an $\operatorname{AsSQS}^{B}(2 p)$ are non-isomorphic.

Example 2.2.21 (Non-existence). Let $p=13$, then $\llbracket \mathbb{Z}_{p} \rrbracket \backslash\{1\}=\{2,3,4,5,6,7\}$.
Since $2^{-1}=7, \chi=3$, and $\frac{1}{\xi+1}=6$, we have $\llbracket \Lambda_{p} \rrbracket=\{2,4,5\}$. Then $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}^{*}\right)$ is a $2-(3,2,2)$ design, where

$$
\begin{aligned}
\mathcal{B}= & \{\{2,3,4\},\{2,3,6\},\{2,5,6\},\{2,5,7\},\{3,4,5\},\{3,5,6\},\{4,6,7\}\} \\
& \cup\{\{2,4\},\{3,7\},\{4,5\},\{4,6\}\}, \\
\mathcal{B}^{*}= & \{\{2,4\},\{2,4\},\{2,5\},\{2,5\},\{4,5\},\{4,5\}\} \\
\mathcal{B}^{\triangle}= & \emptyset
\end{aligned}
$$

Hence, there does not exist an $\operatorname{AsSQS}^{B}(26)$.
Example 2.2.22. Let $p=17$, then $\llbracket \mathbb{Z}_{p} \rrbracket \backslash\{1\}=\{2,3, \ldots, 9\}$. Since $2^{-1}=9$ and $\chi=7$, we have $\llbracket \Lambda_{p} \rrbracket=\{2,3,4,5,6,8\}$. Then $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}^{*}\right)$ is a $(6,\{2,3\}, 2)-\mathrm{PBD}$, where

$$
\begin{aligned}
\mathcal{B}^{\triangle}= & \{\{2,3,4\},\{2,3,8\},\{2,4,5\},\{3,5,6\},\{4,6,8\},\{5,6,8\}\} \\
\mathcal{B}^{*}= & \{\{2,5\},\{2,6\},\{2,6\},\{2,8\},\{3,4\},\{3,5\},\{3,6\},\{3,8\},\{4,5\}, \\
& \{4,6\},\{4,8\},\{5,8\}\} \cup \mathcal{B}^{\triangle}
\end{aligned}
$$

The hypergraph $\left(\llbracket \Lambda_{p} \rrbracket, \mathcal{B}_{2}^{\triangle}, \gamma\right)$ has a rainbow 1-factor

$$
\mathcal{F}=\{\{2,3,4\},\{5,6,8\}\} \subset \mathcal{B}^{\triangle}
$$

such that $\gamma(\{2,3,4\})=C(3 / 2)=C(10)$ and $\gamma(\{5,6,8\})=C(5 / 8)=C(7)$, where $C(10) \cup C(7)=\{3,6,12,10,8,15\} \cup\{4,13,5,7,11,14\}=\Omega_{p}^{*}$. Hence $\{0,1,10 ; 3\}$ and $\{0,1,7 ; 5\}$ can be chosen as Type II base blocks of an $\operatorname{AsSQS}^{B}(34)$.

We have verified the following results on rainbow 1-factors by using computers:

Theorem 2.2.23. An $A s S Q S^{B}(2 p)$ having no Type III quadruples exists if $p \equiv 1,5(\bmod 12)$ is prime and $17 \leq p<1000$.

Actually, the Construction 2.2 .20 can be naturally generalized to "affineinvariant" SQSs over $\mathbb{F}_{q} \oplus \mathbb{F}_{2}$ with a non-prime $q$. We give an example for $q=7^{2}$. This provides a 2-chromatic SQS(98) whose existence is previously unknown (see Ji [59]).

Example 2.2.24. Let $\alpha$ be a primitive element with $\alpha^{2}+1=0$ in $\mathbb{F}_{49}$. Let

$$
\begin{aligned}
& B_{I}=\{0,1 ; 3 \alpha+4,4 \alpha+4\} \\
& B_{0}=\{0,1,6 ; 0\}, B_{\xi}=\{0,1,3 ; 6\}, \text { and } B_{i}=\left\{0,1, a_{i} ; b_{i}\right\}
\end{aligned}
$$

where

$$
\begin{aligned}
\left(a_{1}, a_{2}, \ldots, a_{7}\right) & =(6 \alpha+1, \alpha+4,2 \alpha+2,5 \alpha+3,6 \alpha+6, \alpha+1,5) \text { and } \\
\left(b_{1}, b_{2}, \ldots, b_{7}\right) & =(2 \alpha+3,5 \alpha+5,5 \alpha+2,3 \alpha+6,2 \alpha+1,4 \alpha+5,3 \alpha+1)
\end{aligned}
$$

Let $\mathcal{B}=\left\{m B+c \mid B \in\left\{B_{I}, B_{\xi}, B_{0}, B_{1}, \ldots, B_{7}\right\}, m \in \mathbb{F}_{49}^{\times}, c \in \mathbb{F}_{49}\right\}$, where $\left|\mathcal{O}_{A}\left(B_{I}\right)\right|=12 \times 98,\left|\mathcal{O}_{A}\left(B_{0}\right)\right|=24 \times 98,\left|\mathcal{O}_{A}\left(B_{\xi}\right)\right|=16 \times 98$, and $\left|\mathcal{O}_{A}\left(B_{i}\right)\right|=$ $48 \times 98$ for $1 \leq i \leq 7$. Then, $\left(\mathbb{F}_{49} \oplus \mathbb{F}_{2}, \mathcal{B}\right)$ is a 2-chromatic $\operatorname{SQS}(98)$.

Note that, this $\operatorname{SQS}(98)$ is not cyclic, hence it is not an $\operatorname{AsSQS}(98)$. We will show the non-existence of an $\operatorname{AsSQS}(98)$ in Section 2.4 .

### 2.3 Affine-invariant strictly cyclic Steiner quadruple systems over $\mathbb{Z}_{2 p^{m}}$

### 2.3.1 Preliminaries

For the recursive constructions, we still suppose $p \equiv 1,5(\bmod 12)$ is a prime. Let $m$ be a positive integer. We aim to construct an $\operatorname{AsSQS}\left(2 p^{m}\right)$ based on an $\operatorname{AsSQS}^{A}(2 p)$ or an $\operatorname{AsSQS}^{B}(2 p)$ proposed in Sections 2.2.2 and 2.2.3.

Let $\chi_{m}$ denote a root of $2 \chi_{m}^{2}-2 \chi_{m}+1=0$ over $\mathbb{Z}_{p^{m}}$. For $p \equiv 1(\bmod 12)$, let $\xi_{m}$ denote a root of $\xi_{m}^{2}-\xi_{m}+1=0$ over $\mathbb{Z}_{p^{m}}$. It is easily seen that $\chi_{m} \equiv \chi_{1}$ $(\bmod p)$ and $\xi_{m} \equiv \xi_{1}(\bmod p)$ hold for any positive integer $m$.

Note that the multiplicative group $\mathbb{Z}_{2 p^{m}}^{\times} \cong \mathbb{Z}_{p^{m}}^{\times} \times \mathbb{Z}_{2}^{\times} \cong \mathbb{Z}_{p^{m}}^{\times}$is of order $\varphi\left(p^{m}\right)=p\left(p^{m-1}-1\right)$, where $\mathbb{Z}_{p^{m}}^{\times}=\mathbb{Z}_{p^{m}} \backslash p \mathbb{Z}_{p^{m-1}}=\left(\mathbb{Z}_{p} \backslash\{0\}\right)+p \mathbb{Z}_{p^{m-1}}$ and

$$
p \mathbb{Z}_{p^{m-1}}=p \mathbb{Z} / p^{m} \mathbb{Z}=\left\{p, 2 p, \ldots, p\left(p^{m-1}-1\right)\right\} \quad\left(\bmod p^{m}\right)
$$

Definition 2.3.1. For $t \in[m-1]$, let $h=\min \{t, m-t\}$. Define a group homomorphism $\psi_{t}: \mathbb{Z}_{p^{m-t}}^{\times} \rightarrow \mathbb{Z}_{p^{h}}^{\times}$for any $x \in \mathbb{Z}_{p^{m-t}}^{\times}$, so that

$$
\psi_{t}(x) \equiv x \quad\left(\bmod p^{h}\right)
$$

Conversely, for any $y \in \mathbb{Z}_{p^{h}}^{\times}$, define

$$
\psi_{t}^{-1}(y)=\left\{x \in \mathbb{Z}_{p^{m-t}}^{\times} \mid \psi_{t}(x)=y\right\} .
$$

In particular, $\psi_{t}$ is trivial if $m-t \leq t$.
By means of the above homomorphism $\psi_{t}$, we define two specific subsets of $\mathbb{Z}_{p^{m-t}}^{\times}$, namely $S_{t}$ and $R_{t}$, which provide important parameters in the recursive constructions below. In the rest of this subsection, we sum up a series of properties of $S_{t}$ and $R_{t}$ which allow us to get some partitions related to the multiplicative group $\mathbb{Z}_{p^{m-t}}^{\times}$.

Proposition 2.3.2. Let $g$ be a generator of $\mathbb{Z}_{p^{h}}^{\times}$, where $h=\min \{t, m-t\}$. Let

$$
\begin{equation*}
S_{t}=\bigcup_{k=0}^{\frac{\varphi\left(p^{h}\right)}{2}-1} \psi_{t}^{-1}\left(g^{k}\right) \subset \mathbb{Z}_{p^{m-t}}^{\times} \tag{2.23}
\end{equation*}
$$

where $\psi_{t}$ is defined as in Definition 2.3.1. Then,
(i) $\left|S_{t}\right|=\frac{\varphi\left(p^{m-t}\right)}{2}$.
(ii) $S_{t}+p^{t} \equiv S_{t}\left(\bmod p^{m-t}\right)$.
(iii) $S_{t} \cup\left(-S_{t}\right)=\mathbb{Z}_{p^{m-t}}^{\times}$.

Proof. First, we suppose $m-t \leq t$.
(i) This is straightforward by the definition.
(ii) This follows from $p^{t} \equiv 0\left(\bmod p^{m-t}\right)$ under the assumption $t \geq m-t$.
(iii) Note that $-S_{t}=\left\{\left.g_{t}^{k+\frac{\varphi\left(p^{m-t}\right)}{2}} \right\rvert\, k \in\left[0, \frac{\varphi\left(p^{m-t}\right)}{2}-1\right]\right\}=\left\{g_{t}^{k} \mid k \in\right.$ $\left.\left[\frac{\varphi\left(p^{m-t}\right)}{2}, \varphi\left(p^{m-t}\right)-1\right]\right\}$. Hence $S_{t} \cap\left(-S_{t}\right)=\emptyset$. By combining with (ii), the conclusion is straightforward.

Next, we suppose $m-t>t$.
(i) Note that $\psi_{t}$ is a group homomorphism, hence $\left|\psi_{t}^{-1}(y)\right|=p^{m-2 t}$. A direct calculation gives that $\left|S_{t}\right|=\frac{\varphi\left(p^{t}\right)}{2} p^{m-2 t}=\frac{\varphi\left(p^{m-t}\right)}{2}$.
(ii) This follows from $\psi_{t}^{-1}(y)+p^{t}=\psi_{t}^{-1}(y)$ for any $y \in \mathbb{Z}_{p^{t}}^{\times}$.
(iii) Since $\psi_{t}$ is a group homomorphism, we have

$$
-S_{t}=\bigcup_{k=0}^{\frac{\varphi\left(p^{t}\right)}{2}-1} \psi_{t}^{-1}\left(-g_{m-t}^{k}\right)=\bigcup_{k=\frac{\varphi\left(p^{t}\right)}{2}}^{\varphi\left(p^{t}\right)-1} \psi_{t}^{-1}\left(g_{m-t}^{k}\right)
$$

which implies $S_{t} \cup\left(-S_{t}\right)=\psi_{t}^{-1}\left(\mathbb{Z}_{p^{t}}^{\times}\right)=\mathbb{Z}_{p^{m-t}}^{\times}$.
Example 2.3.3. Let $p=5, m=3$, and $t=1$. Then $S_{1}=\psi^{-1}(1) \cup \psi^{-1}(2)=$ $\{1,6,11,16,21\} \cup\{2,7,12,17,22\} \subset \mathbb{Z}_{5^{2}}^{\times}$.

Next, we introduce a function on $\mathbb{Z}_{p^{m-t}}^{\times}$for our constructions (see Lemma 2.3.19 of Type IV base blocks).

Proposition 2.3.4. For $s \in \mathbb{Z}_{p^{m-t}}^{\times}$, let $\zeta(s)=s+c\left(2 s-p^{t}\right)^{-1} s p^{t}$, where $c$ is a constant in $\mathbb{Z}_{p^{m-t}}^{\times}$. Then $\left\{\zeta(s) \mid s \in S_{t}\right\}=S_{t}$.

Proof. If $m-t \leq t$, then $\left(2 s-p^{t}\right)^{-1} \equiv(2 s)^{-1}\left(\bmod p^{m-t}\right)$. Thus $\zeta(s)=$ $s+2^{-1} c p^{t}$. It follows from Proposition 2.3 .2 (iii) that $\left\{\zeta(s) \mid s \in S_{t}\right\}=S_{t}+$ $2^{-1} c p^{t}=S_{t}$.

If $m-t>t$, then $\psi(\zeta(s))=\psi(s)$ holds for every $s \in \mathbb{Z}_{p^{m-t}}^{\times}$. In words, $\zeta(s) \in \psi^{-1}(\psi(s)) \subseteq S_{t}$ holds for every $s \in S_{t}$, which implies $\left\{\zeta(s) \mid s \in S_{t}\right\} \subseteq S_{t}$. Then it suffices to show that $\zeta\left(s_{1}\right) \neq \zeta\left(s_{2}\right)$ for any distinct $s_{1}, s_{2} \in S_{t}$.
(a) If $s_{1} \not \equiv s_{2}\left(\bmod p^{t}\right)$, since $\zeta(s) \equiv s\left(\bmod p^{t}\right)$, we have $\zeta\left(s_{1}\right) \not \equiv \zeta\left(s_{2}\right)$ $\left(\bmod p^{t}\right)$. More precisely, $\zeta\left(s_{1}\right) \not \equiv \zeta\left(s_{2}\right)\left(\bmod p^{m-t}\right)$.
(b) If $s_{1} \equiv s_{2}\left(\bmod p^{t}\right)$, then $s_{1}=s_{2}+\ell p^{t}$ for some $\ell \in\left[p^{m-2 t}-1\right]$. Assuming $\zeta\left(s_{1}\right)=\zeta\left(s_{2}\right)$, we have $\ell p^{t}+c\left(2 s_{1}-p^{t}\right)^{-1} s_{1} p^{t} \equiv c\left(2 s_{2}-p^{t}\right)^{-1} s_{2} p^{t}$ $\left(\bmod p^{m-t}\right)$. Thus, $\ell+c\left(2-s_{1}^{-1} p^{t}\right)^{-1} \equiv c\left(2-s_{2}^{-1} p^{t}\right)^{-1}\left(\bmod p^{m-2 t}\right)$, i.e.,

$$
\ell\left(2-s_{1}^{-1} p^{t}\right)\left(2-s_{2}^{-1} p^{t}\right) \equiv c\left(s_{2}^{-1}-s_{1}^{-1}\right) p^{t} \quad\left(\bmod p^{m-2 t}\right)
$$

Moreover, note that $s_{2}^{-1}-s_{1}^{-1}=\left(s_{1}-\ell p^{t}\right)^{-1}-s_{1}^{-1}=\ell \sum_{k=1}^{\infty} s_{1}^{-k-1} \ell^{k-1} p^{k t}$, We finally derive

$$
\left(2-s_{1}^{-1} p^{t}\right)\left(2-s_{2}^{-1} p^{t}\right) \equiv c p^{2 t} \sum_{k=0}^{\infty} s_{1}^{-k-2} \ell^{k} p^{k t} \quad\left(\bmod p^{m-2 t-\iota}\right)
$$

where $\iota$ is the non-negative integer satisfying $p^{\iota}=\operatorname{gcd}\left(\ell, p^{m-2 t}\right)$. It is clear that the left-hand side is invertible but the right-hand side is not, which leads to a contradiction. Therefore, $\zeta\left(s_{1}\right) \not \equiv \zeta\left(s_{2}\right)\left(\bmod p^{m-t}\right)$.

Let $H_{0}^{(t)}$ denote the multiplicative subgroup of $\mathbb{Z}_{p^{m}}^{\times}$generated by $g_{0}^{\varphi\left(p^{m-t}\right)}$. Then, we denote the cosets of $H_{0}^{(t)}$ by

$$
\begin{equation*}
H_{a}^{(t)}=g_{0}^{a} H_{0}^{(t)}=\left\{g_{0}^{a+k \varphi\left(p^{m-t}\right)} \mid k \in\left[0, p^{t}-1\right]\right\} \tag{2.24}
\end{equation*}
$$

Proposition 2.3.5. For any $x, y \in H_{0}^{(t)}$ and any $z \in H_{\frac{\varphi\left(p^{m-t}\right)}{2}}^{(t)}$, the following hold:
(i) $x-y \equiv 0\left(\bmod p^{m-t}\right)$.
(ii) $x+z \equiv 0\left(\bmod p^{m-t}\right)$.

Proof. Suppose $x=g_{0}^{k_{1} \varphi\left(p^{m-t}\right)}, y=g_{0}^{k_{2} \varphi\left(p^{m-t}\right)}$, and $z=g_{0}^{k_{3} \varphi\left(p^{m-t}\right)+\frac{\varphi\left(p^{m-t}\right)}{2}}$. Since $g_{0}^{\varphi\left(p^{m-t}\right)} \equiv 1\left(\bmod p^{m-t}\right)$, we have

$$
g_{0}^{\varphi\left(p^{m-t}\right)} p^{t} \equiv p^{t} \quad\left(\bmod p^{m}\right) \quad \text { and } \quad g_{0}^{\frac{\varphi\left(p^{m-t}\right)}{2}} p^{t} \equiv-p^{t} \quad\left(\bmod p^{m}\right)
$$

Then,
(i) $x p^{t}-y p^{t} \equiv p^{t}-p^{t} \equiv 0\left(\bmod p^{m}\right)$, which implies $x-y \equiv 0\left(\bmod p^{m-t}\right)$;
(ii) $x p^{t}+z p^{t} \equiv p^{t}-p^{t} \equiv 0\left(\bmod p^{m}\right)$, which implies $x+z \equiv 0\left(\bmod p^{m-t}\right)$.

Proposition 2.3.6. The following holds:

$$
\left(\bigcup_{s \in S_{t}} s H_{0}^{(t)}\right) \bigcup\left(\bigcup_{s \in S_{t}} s H_{\frac{\varphi\left(p^{m-t}\right)}{2}}^{(t)}\right)=\mathbb{Z}_{p^{m}}^{\times}
$$

Proof. Note that the left-hand side is a union of cosets, in which the number of cosets is $2\left|S_{t}\right|=\varphi\left(p^{m-t}\right)$ and each coset is of size $p^{t}$. Hence it suffices to show those cosets are mutually disjoint. Clearly, for some given $s \in S_{t}, s H_{0}^{(t)}$ and $s H_{\frac{\varphi\left(p^{m-t}\right)}{2}}^{(t)}$ never coincide. Now we give the proofs for the remaining cases by contradiction.
(a) For distinct $s_{1}, s_{2} \in S_{t}$, assume $s_{1} H_{0}^{(t)}=s_{2} H_{0}^{(t)}$. There must exist $x, y \in$ $H_{0}^{(t)}$, such that $s_{1} x \equiv s_{2} y\left(\bmod p^{m}\right)$. It is known by Proposition 2.3.5 (i) that there exists $r \neq 0$, such that $y=x+r p^{m-t}$. Hence, we have $s_{1} x \equiv$ $s_{2}\left(x+r p^{m-t}\right)\left(\bmod p^{m}\right)$, i.e., $\left(s_{1}-s_{2}\right) x \equiv s_{2} r p^{m-t}\left(\bmod p^{m}\right)$. Since $x$ is invertible, there must be $s_{1}-s_{2} \equiv 0\left(\bmod p^{m-t}\right)$. This is known to be impossible because $s_{1}, s_{2}$ are distinct in $\mathbb{Z}_{p^{m-t}}^{\times}$. In the same manner, it can be easily shown that $s_{1} H_{a}^{(t)} \neq s_{2} H_{a}^{(t)}$ for any given $a$ and distinct $s_{1}, s_{2} \in S_{t}$.
(b) For distinct $s_{1}, s_{2} \in S_{t}$, assume $s_{1} H_{0}^{(t)}=s_{2} H_{\frac{\varphi\left(p^{m-t}\right)}{2}}^{(t)}$. Then there exist $x \in H_{0}^{(t)}$ and $z \in H_{\frac{\varphi\left(p^{m-t}\right)}{2}}^{(t)}$ so that $s_{1} x \equiv s_{2} z\left(\bmod p^{m}\right)$. By Proposition 2.3.5 (iii), we can suppose $z=r p^{m-t}-x$ for some $r$. Then we have $\left(s_{1}+s_{2}\right) x \equiv s_{2} r p^{m-t}\left(\bmod p^{m}\right)$ which implies $s_{1}+s_{2} \equiv 0\left(\bmod p^{m-t}\right)$. On one hand, $s_{1} \in S_{t}$ requires $s_{2}=-s_{1} \in-S_{t}$, but on the other hand, $s_{2} \in S_{t}$. It is known by Proposition 2.3 .2 that $S_{t} \cap\left(-S_{t}\right)=\emptyset$, hence $s_{1}+s_{2} \equiv 0$ $\left(\bmod p^{m-t}\right)$ cannot hold for any distinct $s_{1}, s_{2} \in S_{t}$.

For $p \equiv 1(\bmod 12)$, we introduce another subset of $\mathbb{Z}_{p^{m-t}}^{\times}$, which consists of one third of the elements of $\mathbb{Z}_{p^{m-t}}^{\times}$.

Proposition 2.3.7. Let $g$ be a generator of $\mathbb{Z}_{p^{h}}^{\times}$, where $h=\min \{t, m-t\}$. Let

$$
\begin{equation*}
R_{t}=\left(\bigcup_{k=0}^{\frac{\varphi\left(p^{h}\right)}{6}-1} \psi_{t}^{-1}\left(g^{k}\right)\right) \bigcup\left(\bigcup_{k=\frac{\varphi\left(p^{h}\right)}{2}}^{\frac{2 \varphi\left(p^{h}\right)}{3}-1} \psi_{t}^{-1}\left(g^{k}\right)\right) \subset \mathbb{Z}_{p^{m-t}}^{\times} \tag{2.25}
\end{equation*}
$$

where $\psi_{t}$ is defined as in Definition 2.3.1. Then,
(i) $\left|R_{t}\right|=\frac{\varphi\left(p^{m-t}\right)}{3}$.
(ii) $R_{t}+p^{t} \equiv R_{t}\left(\bmod p^{m-t}\right)$.
(iii) $R_{t} \cup \xi_{m}^{2} R_{t} \cup \xi_{m}^{4} R_{t}=\mathbb{Z}_{p^{m-t}}^{\times}$.

Proof. We omit the proofs of (ii) and (ii), because they are the same as the proofs of Proposition 2.3.2. By recalling that $\xi_{m}^{2} \equiv g_{t}^{\frac{\varphi\left(p^{m-t}\right)}{3}}\left(\bmod p^{m-t}\right)$ and $\xi_{m}^{2} \equiv g_{m}^{\frac{\varphi\left(p^{t}\right)}{-t}}\left(\bmod p^{t}\right)$, we prove (iii) in the following two cases:

For the first case, suppose $m-t \leq t$. Let

$$
I=\left[0, \frac{\varphi\left(p^{m-t}\right)}{6}-1\right] \cup\left[\frac{\varphi\left(p^{m-t}\right)}{2}, \frac{2 \varphi\left(p^{m-t}\right)}{3}-1\right]
$$

denote a subset of $\mathbb{Z}_{\varphi\left(p^{m-t}\right)}$, then $I \cup\left(I+\frac{\varphi\left(p^{m-t}\right)}{3}\right) \cup\left(I+\frac{2 \varphi\left(p^{m-t}\right)}{3}\right)=\mathbb{Z}_{\varphi\left(p^{m-t}\right)}$, which implies $R_{t} \cup \xi_{m}^{2} R_{t} \cup \xi_{m}^{4} R_{t}=\left\{g_{t}^{k} \mid k \in \mathbb{Z}_{\varphi\left(p^{m-t}\right)}\right\}=\mathbb{Z}_{p^{m-t}}^{\times}$.

For the second case, suppose $m-t>t$. Let

$$
T=\left\{g_{m-t}^{k} \left\lvert\, k \in\left[0, \frac{\varphi\left(p^{t}\right)}{6}-1\right] \cup\left[\frac{\varphi\left(p^{t}\right)}{2}, \frac{2 \varphi\left(p^{t}\right)}{3}-1\right]\right.\right\} .
$$

In the same manner as the first case, we have $T \cup \xi_{m}^{2} T \cup \xi_{m}^{4} T=\mathbb{Z}_{p^{t}}^{\times}$. Moreover, since $\psi$ is a group homomorphism, we obtain $R_{t} \cup \xi_{m}^{2} R_{t} \cup \xi_{m}^{4} R_{t}=\psi^{-1}\left(\mathbb{Z}_{p^{t}}^{\times}\right)=$ $\mathbb{Z}_{p^{m-t}}^{\times}$. By combining with (i), it is easy to see they are mutually disjoint.

Next, we define a function on $\mathbb{Z}_{p^{m-t}}^{\times}$for our constructions when $p \equiv 1$ (mod 12) (see Lemma 2.3.20 for Type III ${ }^{\xi}$ base blocks and Lemma 2.3.24 for Type $\mathrm{II}^{\xi}$ base blocks). The coefficient $\sqrt{-3}$ is considered as an element in $\mathbb{Z}_{p^{m-t}}$ whose square is -3 . Note that there are two elements satisfying the above condition, hence we should choose $\sqrt{-3}$ as in $\xi_{m}=\frac{1+\sqrt{-3}}{2}$. Without loss of generality, here we set $\sqrt{-3}=2 \xi_{m}-1$.

Proposition 2.3.8. Let $\vartheta(s)=\left(3-\sqrt{-3} s p^{t}\right)^{-1} s$. Then

$$
\bigcup_{s \in R_{t}}\left\{\vartheta(s), \xi_{m}^{2} \vartheta(s), \xi_{m}^{4} \vartheta(s)\right\}=\mathbb{Z}_{p^{m-t}}^{\times}
$$

Proof. Let $R_{t}^{*}=\bigcup_{s \in R_{t}}\left\{s^{-1}\right\}$. It is easy to see that $R_{t}^{*}=\xi_{m}^{2} R_{t}$. By Proposition 2.3.7 (iii), we have $R_{t}^{*}+p^{t}=R_{t}^{*}$. Then,

$$
\bigcup_{s \in R_{t}}\{3 \vartheta(s)\}=\bigcup_{s^{*} \in R_{t}^{*}}\left\{\left(s^{*}-\frac{\sqrt{-3}}{3} p^{t}\right)^{-1}\right\}=\bigcup_{s^{* *} \in R_{t}^{*}}\left\{\left(s^{* *}\right)^{-1}\right\}=R_{t}
$$

Furthermore, it follows from Proposition 2.3 .7 (iii) that

$$
\bigcup_{s \in R_{t}}\left\{\vartheta(s), \xi_{m}^{2} \vartheta(s), \xi_{m}^{4} \vartheta(s)\right\}=3^{-1}\left(R_{t} \cup \xi_{m}^{2} R_{t} \cup \xi_{m}^{4} R_{t}\right)=\mathbb{Z}_{p^{m-t}}^{\times}
$$

Proposition 2.3.9. For any $t \in[m-1]$,

$$
\bigcup_{s \in R_{t}} C\left(\xi_{m}+s p^{t}\right)=\left\{\xi_{m}, 1-\xi_{m}\right\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}
$$

Proof. For any $s \in \mathbb{Z}_{p^{m-t}}^{\times}$, there are three elements in $C\left(\xi_{m}+s p^{t}\right)$ which are congruent to $\xi_{m}$ modulo $p^{t}$, namely, $\xi_{m}+s p^{t}, 1-\left(\xi_{m}+s p^{t}\right)^{-1}$, and $\left(1-\xi_{m}-\right.$ $\left.s p^{t}\right)^{-1}$. Let

$$
\begin{aligned}
& W_{1}=\bigcup_{s \in R_{t}}\left\{\xi_{m}+s p^{t}\right\}, \\
& W_{2}=\bigcup_{s \in R_{t}}\left\{1-\left(\xi_{m}+s p^{t}\right)^{-1}\right\}, \text { and } \\
& W_{3}=\bigcup_{s \in R_{t}}\left\{\left(1-\xi_{m}-s p^{t}\right)^{-1}\right\}
\end{aligned}
$$

be three subsets of $\xi_{m}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}$with the same size $\frac{\varphi\left(p^{m-t}\right)}{3}$. Now we proof their disjointness by contradiction. Assume there exist $s_{1}, s_{2}, s_{3} \in R_{t}$, such that at least one of the following holds:

$$
\begin{aligned}
& \xi_{m}+s_{1} p^{t}=1-\left(\xi_{m}+s_{2} p^{t}\right)^{-1} \\
& 1-\left(\xi_{m}+s_{2} p^{t}\right)^{-1}=\left(1-\xi_{m}-s_{3} p^{t}\right)^{-1} \\
& \left(1-\xi_{m}-s_{3} p^{t}\right)^{-1}=\xi_{m}+s_{1} p^{t}
\end{aligned}
$$

which are equivalent to

$$
\begin{align*}
& s_{1} \equiv \xi_{m}^{4} s_{2}+\xi_{m}^{2} s_{1} s_{2} p^{t} \quad\left(\bmod p^{m-t}\right)  \tag{2.26}\\
& s_{2} \equiv \xi_{m}^{4} s_{3}+\xi_{m}^{2} s_{2} s_{3} p^{t} \quad\left(\bmod p^{m-t}\right)  \tag{2.27}\\
& s_{3} \equiv \xi_{m}^{4} s_{1}+\xi_{m}^{2} s_{3} s_{1} p^{t} \quad\left(\bmod p^{m-t}\right) \tag{2.28}
\end{align*}
$$

Without loss of generality, we assume 2.26 holds. Then we have

$$
\begin{equation*}
\psi_{t}\left(s_{1}\right) \equiv \psi_{t}\left(\xi_{m}^{4}\right) \psi_{t}\left(s_{2}\right) \quad\left(\bmod p^{h}\right) \tag{2.29}
\end{equation*}
$$

where $h=\min \{t, m-t\}$. Here $\psi_{t}\left(\xi_{m}^{4}\right)$ can be regarded as the $\frac{2 \varphi\left(p^{h}\right)}{3}$-th power of the generator of $\mathbb{Z}_{p^{h}}^{\times}$. However, for $I=\left[0, \frac{\varphi\left(p^{h}\right)}{6}-1\right] \cup\left[\frac{\varphi\left(p^{h}\right)}{2}, \frac{2 \varphi\left(p^{h}\right)}{3}-1\right] \subset \mathbb{Z}_{\varphi\left(p^{h}\right)}$, $I \cap\left(I+\frac{2 \varphi\left(p^{h}\right)}{3}\right)=\emptyset$ must hold. Therefore, the congruence 2.29 cannot hold. This proves the disjointness of $W_{1}$ and $W_{2}$. In the same manner, it can be shown that $W_{1}, W_{2}$, and $W_{3}$ are mutually disjoint. Hence, $\left\{W_{1}, W_{2}, W_{3}\right\}$ forms a partition of $\xi_{m}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}$. As a direct consequence, $\left\{1-W_{1}, 1-W_{2}, 1-W_{3}\right\}$
forms a partition of $1-\xi_{m}-p^{t} \mathbb{Z}_{p^{m-t}}^{\times}=\left(1-\xi_{m}\right)+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}$, where

$$
\begin{aligned}
& 1-W_{1}=\bigcup_{s \in R_{t}}\left\{1-\xi_{m}-s p^{t}\right\} \\
& 1-W_{2}=\bigcup_{s \in R_{t}}\left\{\left(\xi_{m}+s p^{t}\right)^{-1}\right\} \\
& 1-W_{3}=\bigcup_{s \in R_{t}}\left\{1-\left(1-\xi_{m}-s p^{t}\right)^{-1}\right\}
\end{aligned}
$$

which are also subsets of $\bigcup_{s \in R_{t}} C\left(\xi_{m}+s p^{t}\right)$.

### 2.3.2 Recursive construction A

In this subsection, we provide a recursive construction for an $\operatorname{AsSQS}\left(2 p^{m}\right)$ based on $\operatorname{AsSQS}^{A}(2 p)$. We denote the resulting sSQS by $\operatorname{AsSQS}^{A}\left(2 p^{m}\right)$.

Construction 2.3.10. Assume that both an $\operatorname{AsSQS}^{A}(2 p)$ and an $\operatorname{AsSQS}^{A}\left(2 p^{m-1}\right)$ have been constructed, then the base blocks of an $\operatorname{AsSQS}^{A}\left(2 p^{m}\right)$ can be obtained as follows:
(i) For the cases of prime $p \equiv 1$ or $5(\bmod 12)$, we have the base blocks as follows:

Type I': $\left\{0,1 ; \chi_{m}, 1-\chi_{m}\right\}$;
Type I: $\left\{0,1 ; b_{i}+s p^{m-1}, 1-\left(b_{i}+s p^{m-1}\right)\right\}$ for $i \in\left[\frac{p-5}{4}\right]$ and $s \in[0, p-1]$;
Type II': $\left\{0,1,-1+s p^{m-1} ; s p^{m-1}\right\}$ for $s \in\left[0, \frac{p-1}{2}\right]$;
Type IV: $\left\{0, p^{t}, s_{t} ; \chi_{m} s_{t}+\left(2 s_{t}-p^{t}\right)^{-1}\left(1-\chi_{m}\right) p^{t} s_{t}\right\}$ for $t \in[m-1]$ and $s_{t} \in S_{t}$, where $S_{t}$ is defined by 2.23 ;
Type V: $p B\left(\bmod p^{m}\right)$ for every base block $B$ of an $\operatorname{AsSQS}\left(2 p^{m-1}\right)$.
(ii) If $p \equiv 5(\bmod 12)$, we additionally have

Type III: $\left\{0,1, a_{i}+s p^{m-1}, 1-\left(a_{i}+s p^{m-1}\right)\right\}$ for $i \in\left[\frac{p-5}{12}\right]$ and $s \in[0, p-1]$;
(iii) If $p \equiv 1(\bmod 12)$, we additionally have

Type III: $\left\{0,1, a_{i}+s p^{m-1}, 1-\left(a_{i}+s p^{m-1}\right)\right\}$ for $i \in\left[\frac{p-13}{12}\right]$ and $s \in$ $[0, p-1]$;
Type $\operatorname{III}{ }^{\xi}:\left\{0,1, \xi_{m}, \overline{\xi_{m}}\right\}$ and $\left\{0,1, \xi_{m}+s_{t} p^{t}, \overline{\xi_{m}}+\left(3-\sqrt{-3} s_{t} p^{t}\right)^{-1} s_{t} p^{t}\right\}$ for $s_{t} \in R_{t}$, where $R_{t}$ is defined by 2.25 .

Example 2.3.11. Let $p=5$ and $m=2$. Take $\chi_{m}=4$. The base blocks are given as follows:

Type I': $\{0,1 ; 4,22\}$,
Type $\mathrm{II}^{\prime}:\{0,1,24 ; 0\},\{0,1,4 ; 5\},\{0,1,9 ; 10\}$,

Type IV: $\{0,5,1 ; 9\},\{0,5,2 ; 13\}$,
Type V: $\{0,5 ; 10,20\},\{0,5,20 ; 0\}$.
Example 2.3.12. Let $p=5$ and $m=3$. Take $\chi_{m}=29$. The base blocks are given as follows:

Type I': $\{0,1 ; 29,97\}$,
Type II': $\{0,1,124 ; 0\},\{0,1,24 ; 25\},\{0,1,49 ; 50\}$,
Type IV: $\{0,5,1 ; 34\},\{0,5,6 ; 54\},\{0,5,11 ; 74\},\{0,5,16 ; 94\},\{0,5,21 ; 114\}$, $\{0,5,2 ; 88\},\{0,5,7 ; 108\},\{0,5,12 ; 3\},\{0,5,17 ; 23\},\{0,5,22 ; 43\}$, $\{0,25,1 ; 54\},\{0,25,2 ; 83\}$,

Type V: $5 B(\bmod 125)$ for every $B \in \mathcal{B}$, where $\mathcal{B}$ consists of all base blocks listed in Example 2.3.11.

It is clear that every pure triple is a member of $\mathcal{O}_{A}(\{0,1, a\})$ for some $a$, and every mixed triple is a member of $\mathcal{O}_{A}(\{0,1 ; b\})$ for some $b$. Now we begin to prove that every pure or mixed triple can be covered exactly once in $\mathcal{O}_{A}(B)$ for some base block $B$ in Construction 2.3.10.

Lemma 2.3.13 (Type $\left.\mathrm{I}^{\prime}\right)$. There are exactly two mixed triples containing $\{0,1\}$ covered by $\mathcal{O}_{A}\left(\left\{0,1 ; \chi_{m}, 1-\chi_{m}\right\}\right)$, namely, $\left\{0,1 ; \chi_{m}\right\}$ and $\left\{0,1 ; 1-\chi_{m}\right\}$.

Proof. Note that $\chi_{m}$ is a root of $2 \chi_{m}^{2}-2 \chi_{m}+1=0$ over $\mathbb{Z}_{p^{m}}$, hence $\overline{\chi_{m}}=\chi_{m}$. Therefore, $\operatorname{orb}_{\boldsymbol{A C}}\left(\chi_{m}\right)=\left\{\chi_{m}, 1-\chi_{m}\right\}$ and $\mathcal{O}_{A}\left(\left\{0,1 ; \chi_{m}, 1-\chi_{m}\right\}\right)$ covers only two mixed triples containing $\{0,1\}$.

Lemma 2.3.14 (Type I). Let $B_{1}\left(b_{i}, s\right)=\left\{0,1 ; b_{i}+s p^{m-1}, 1-\left(b_{i}+s p^{m-1}\right)\right\}$. Then, $\bigcup_{i=1}^{\frac{p-5}{4}} \bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{1}\left(b_{i}, s\right)\right)$ covers each mixed triple in

$$
\left\{\{0,1 ; y\} \mid y \in\left(\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi_{m}, 1-\chi_{m}\right\}\right)+p \mathbb{Z}_{p^{m-1}}\right\}
$$

exactly once.
Proof. Given $b_{i}$ and $s, \mathcal{O}_{A}\left(B_{1}\left(b_{i}, s\right)\right)$ covers four mixed triples containing $\{0,1\}$ exactly once, that is, $\left\{\{0,1 ; y\} \mid y \in \operatorname{orb}_{\boldsymbol{A C}}\left(b_{i}+s p^{m-1}\right)\right\}$. Thus, it suffices to prove the disjointness of $\operatorname{orb}_{\boldsymbol{A C}}\left(b_{i_{1}}+s_{1} p^{m-1}\right)$ and $\operatorname{orb}_{\boldsymbol{A C}}\left(b_{i_{2}}+s_{2} p^{m-1}\right)$ for any distinct pairs $\left(i_{1}, s_{1}\right)$ and $\left(i_{2}, s_{2}\right)$. First, when $i_{1} \neq i_{2}$, assume $\operatorname{orb}_{\boldsymbol{A C}}\left(b_{i_{1}}+\right.$ $\left.s p^{m-1}\right)$ and $\operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i_{2}}+s p^{m-1}\right)$ coincide, then $\operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i_{1}}+s p^{m-1}\right) \equiv \operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i_{2}}+\right.$ $\left.s p^{m-1}\right)(\bmod p)$, i.e., $\operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i_{1}}\right) \equiv \operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i_{2}}\right)(\bmod p)$, which contradicts Lemma 2.2.2. Next, when $s_{1} \neq s_{2}$, assume $\operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i}+s_{1} p^{m-1}\right)$ and $\operatorname{orb}_{\boldsymbol{A} \boldsymbol{C}}\left(b_{i}+\right.$ $s_{2} p^{m-1}$ ) coincide, then one of the following must hold:
(a) $b_{i}+s_{1} p^{m-1}=1-\left(b_{i}+s_{2} p^{m-1}\right)$,
(b) $b_{i}+s_{1} p^{m-1}=-\left(b_{i}+s_{2} p^{m-1}\right)\left(1-2\left(b_{i}+s_{2} p^{m-1}\right)\right)^{-1}$,
(c) $b_{i}+s_{1} p^{m-1}=\left(1-\left(b_{i}+s_{2} p^{m-1}\right)\right)\left(1-2\left(b_{i}+s_{2} p^{m-1}\right)\right)^{-1}$,
which implies $b_{i} \equiv 2^{-1}, b_{i} \equiv 0$ or 1 , and $b_{i} \equiv \chi_{1}$ or $1-\chi_{1}(\bmod p)$, respectively. This again contradicts Lemma 2.2.2.

Lemma 2.3.15 (Type $\left.\mathrm{II}^{\prime}\right)$. Let $B_{2}(s)=\left\{0,1,-1+s p^{m-1} ; s p^{m-1}\right\}$ for $s \in$ $\left[0, \frac{p-1}{2}\right]$. Then, $\bigcup_{s=0}^{\frac{p-1}{2}} \mathcal{O}_{A}\left(B_{2}(s)\right)$ covers the pure triple $\{0,1, x\}$ for every $x \in$ $\left\{-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ and the mixed triple $\{0,1 ; y\}$ for every $y \in\left\{0,1,2^{-1}\right\}+$ $p \mathbb{Z}_{p^{m-1}}$ exactly once.

Proof. By acting proper affine transformations on $B_{2}(s)$, we obtain all quadruples containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{2}(s)\right)$ as follows:

$$
\begin{aligned}
& \left\{0,1,-1+s p^{m-1} ; s p^{m-1}\right\},\left\{0,-1-s p^{m-1}, 1 ;-s p^{m-1}\right\} \\
& \left\{1,0,2-s p^{m-1} ; 1-s p^{m-1}\right\},\left\{1,2+s p^{m-1}, 0 ; 1+s p^{m-1}\right\} \\
& \left\{\frac{1}{2}-\frac{1}{4} s p^{m-1}, 1,0 ; \frac{1}{2}+\frac{1}{4} s p^{m-1}\right\},\left\{\frac{1}{2}+\frac{1}{4} s p^{m-1}, 0,1 ; \frac{1}{2}-\frac{1}{4} s p^{m-1}\right\} .
\end{aligned}
$$

It is readily checked that the pure triples arising from these base blocks cover $\left\{\{0,1, x\} \mid x \in\left\{-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}\right\}$ exactly once, and the mixed triples cover $\left\{\{0,1 ; y\} \mid y \in\left\{0,1,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}\right\}$ exactly once.

Lemma 2.3.16. Let $\left(x_{1}, s_{1}\right),\left(x_{2}, s_{2}\right)$ be distinct pairs in $\mathbb{Z}_{p} \times[0, p-1]$. Assume $C\left(x_{1}\right), C\left(x_{2}\right) \notin\left\{C(0), C(-1), C\left(\xi_{1}\right)\right\}$. Moreover, assume $\bar{C}\left(x_{1}\right)$ and $\bar{C}\left(x_{2}\right)$ are disjoint if $x_{1} \neq x_{2}$. Then, $\bar{C}\left(x_{1}+s_{1} p^{m-1}\right)$ and $\bar{C}\left(x_{2}+s_{2} p^{m-1}\right)$ are disjoint.

Proof. Firstly, suppose $x_{1} \neq x_{2}$. If $\bar{C}\left(x_{1}+s_{1} p^{m-1}\right)$ and $\bar{C}\left(x_{2}+s_{2} p^{m-1}\right)$ intersect, then $\bar{C}\left(x_{1}+s_{1} p^{m-1}\right)$ and $\bar{C}\left(x_{2}+s_{2} p^{m-1}\right)(\bmod p)$ must intersect, which implies $\bar{C}\left(x_{1}\right) \cap \bar{C}\left(x_{2}\right) \neq \emptyset$ and contradicts the assumption. Secondly, suppose $x_{1}=$ $x_{2}=x$ and $s_{1} \neq s_{2}$. Assume $\bar{C}\left(x+s_{1} p^{m-1}\right)$ and $\bar{C}\left(x+s_{2} p^{m-1}\right)$ intersect. Without loss of generality, assume $C\left(x+s_{1} p^{m-1}\right)=C\left(x+s_{2} p^{m-1}\right)$. Then we can derive that $x$ must be in $\left\{\infty, 2^{-1}, 1,-1, \xi_{1}, 1-\xi_{1}, 0,2\right\}$ which is contrary to the assumption. To sum up, $\bar{C}\left(x_{1}+s_{1} p^{m-1}\right)$ and $\bar{C}\left(x_{2}+s_{2} p^{m-1}\right)$ must be disjoint.

Lemma 2.3.17 (Type III). Let $B_{3}\left(a_{i}, s\right)=\left\{0,1, a_{i}+s p^{m-1}, 1-\left(a_{i}+s p^{m-1}\right)\right\}$ for $s \in[0, p-1]$. Then, $\bigcup_{i=1}^{\ell_{p}} \bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{3}\left(a_{i}, s\right)\right)$ covers every pure triple in

$$
\left\{\{0,1, x\} \mid x \in\left(\mathbb{Z}_{p} \backslash \Theta\right)+p \mathbb{Z}_{p^{m-1}}\right\}
$$

exactly once, where

$$
\begin{array}{r}
\ell_{p}:=\left\{\begin{array}{lll}
\frac{p-5}{12}, & \text { if } p \equiv 5 & (\bmod 12), \\
\frac{p-13}{12}, & \text { if } p \equiv 1 & (\bmod 12),
\end{array}\right. \\
\Theta:=\left\{\begin{array}{lll}
\left\{0,1,-1,2,2^{-1}\right\}, & \text { if } p \equiv 5 & (\bmod 12), \\
\left\{0,1,-1,2,2^{-1}\right\} \cup \bar{C}\left(\xi_{m}\right), & \text { if } p \equiv 1 & (\bmod 12) .
\end{array}\right.
\end{array}
$$

Proof. It is known that $\mathcal{O}_{A}\left(B_{3}\left(a_{i}, s\right)\right)$ covers the pure triple $\{0,1, x\}$ for every $x \in \bar{C}\left(a_{i}+s p^{m-1}\right)$. First, $C\left(a_{i}+s p^{m-1}\right)$ and $C\left(\overline{a_{i}+s p^{m-1}}\right)$ must be disjoint by Lemma 2.2.4. Moreover, note that $a_{i}$ 's satisfy the assumptions in Lemma 2.3.16 by Lemma 2.2.4. Thus, we can conclude the disjointness of $\bar{C}\left(a_{i_{1}}+s_{1} p^{m-1}\right)$ and $\bar{C}\left(a_{i_{2}}+s_{2} p^{m-1}\right)$ for any distinct pairs $\left\{i_{1}, s_{1}\right\}$ and $\left\{i_{2}, s_{2}\right\}$. Therefore, we have

$$
\begin{equation*}
\bigcup_{s=0}^{p-1} \bigcup_{i=1}^{\ell} \bar{C}\left(a_{i}+s p^{m-1}\right) \subseteq\left(\mathbb{Z}_{p} \backslash \Theta\right)+p \mathbb{Z}_{p^{m-1}} \tag{2.30}
\end{equation*}
$$

Furthermore, both the left-hand side and the right-hand side of 2.30 have cardinalities $12 \ell p^{m-1}$, which completes the proof.

Lemma 2.3.18 (Type IV). Let $S_{t}$ be the subset defined in 2.23). Denote $\alpha=\chi_{m}, \beta=1-\chi_{m}$, and $B_{2}^{(t)}\left(s_{t}\right)=\left\{0, p^{t}, s_{t} ; \alpha s_{t}+\left(2 s_{t}-p^{t}\right)^{-1} \beta p^{t} s_{t}\right\}$ for $t \in[m-1]$ and $s_{t} \in S_{t}$. Then, $\bigcup_{t=1}^{m-1} \bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{2}^{(t)}\left(s_{t}\right)\right)$ covers the pure triple $\{0,1, x\}$ for every $\left.x \in\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right\}\right) \backslash\{0,1\}$, and the mixed triple $\{0,1 ; y\}$ for every $\left.y \in\left(\{\alpha, \beta\}+p \mathbb{Z}_{p^{m-1}}\right\}\right) \backslash\{\alpha, \beta\}$ exactly once.

Proof. We exhaustively list all the quadruples containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{2}^{(t)}\left(s_{t}\right)\right)$ as follows:

$$
\begin{aligned}
& B_{2}^{(t)}\left(s_{t}\right) \times s_{t}^{-1}=\left\{0, s_{t}^{-1} p^{t}, 1 ; \alpha+\left(2 s_{t}-p^{t}\right)^{-1} \beta p^{t}\right\} \\
& \left.\left(B_{2}^{(t)}\left(s_{t}\right)\right)-p^{t}\right) \times\left(s_{t}-p^{t}\right)^{-1}=\left\{-\left(s_{t}-p^{t}\right)^{-1} p^{t}, 0,1 ; \alpha-\left(2 s_{t}-p^{t}\right)^{-1} \beta p^{t}\right\} \\
& B_{2}^{(t)}\left(s_{t}\right) \times\left(-s_{t}^{-1}\right)+1=\left\{1,1-s_{t}^{-1} p^{t}, 0 ; \beta-\left(2 s_{t}-p^{t}\right)^{-1} \beta p^{t}\right\} \\
& \left(p^{t}-B_{2}^{(t)}\left(s_{t}\right)\right) \times\left(s_{t}-p^{t}\right)^{-1}+1=\left\{1+\left(s_{t}-p^{t}\right)^{-1} p^{t}, 1,0 ; \beta+\left(2 s_{t}-p^{t}\right)^{-1} \beta p^{t}\right\}
\end{aligned}
$$

Each pure triple in

$$
\left\{\{0,1, x\} \mid x \in\left\{s_{t}^{-1} p^{t},-\left(s_{t}-p^{t}\right)^{-1} p^{t}, 1-s_{t}^{-1} p^{t}, 1+\left(s_{t}-p^{t}\right)^{-1} p^{t}\right\}\right\}
$$

is covered exactly once in $\mathcal{O}_{A}\left(B_{2}^{(t)}\left(s_{t}\right)\right)$. By Proposition 2.3.2,

$$
\left\{s_{t}^{-1},-\left(s_{t}-p^{t}\right)^{-1} \quad\left(\bmod p^{m-t}\right) \mid s_{t} \in S_{t}\right\}=S_{t} \cup\left(-S_{t}\right)=\mathbb{Z}_{p^{m-t}}^{\times}
$$

Therefore, for each $t \in[m-1]$ and each $s_{t} \in S_{t}$, the pure triple $\{0,1, x\}$ is covered exactly once for every

$$
\left.x \in \bigcup_{t=1}^{m-1}\left(\{0,1\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}\right)=\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right\}\right) \backslash\{0,1\}
$$

Each mixed triple in

$$
\left\{\{0,1 ; y\} \mid y \in\left\{\alpha \pm\left(2 s_{t}-p^{t}\right)^{-1} \beta p^{t}, \beta \pm\left(2 s_{t}-p^{t}\right)^{-1} \beta p^{t}\right\}\right\}
$$

is covered exactly once in $\mathcal{O}_{A}\left(B_{2}^{(t)}\left(s_{t}\right)\right)$. By Proposition 2.3.2,

$$
\left\{ \pm\left(2 s_{t}-p^{t}\right)^{-1} \quad\left(\bmod p^{m-t}\right) \mid s_{t} \in S_{t}\right\}=\mathbb{Z}_{p^{m-t}}^{\times}
$$

Therefore, for each $t \in[m-1]$ and each $s_{t} \in S_{t}$, the mixed triple $\{0,1 ; y\}$ is covered exactly once for every

$$
\left.y \in \bigcup_{t=1}^{m-1}\left(\{\alpha, \beta\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}\right)=\left(\{\alpha, \beta\}+p \mathbb{Z}_{p^{m-1}}\right\}\right) \backslash\{\alpha, \beta\}
$$

Lemma 2.3.19 (Type IV). For a given $t \in[m-1], \bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{2}^{(t)}\left(s_{t}\right)\right)$ covers the pure triple $\left\{0, p^{t}, x_{t}\right\}$ for every $x_{t} \in \mathbb{Z}_{p^{m}}^{\times}$, and the mixed triple $\left\{0, p^{t} ; y_{t}\right\}$ for every $y_{t} \in \mathbb{Z}_{p^{m}}^{\times}$exactly once.

Proof. Let $g_{0}$ be a generator of $\mathbb{Z}_{p^{m}}^{\times}$. We denote $q=p^{m-t}$. Let

$$
\begin{aligned}
Q_{1}(s, u) & =B_{2}^{(t)}(s) \times g_{0}^{u \varphi(q)}, \\
& =\left\{0, p^{t}, g_{0}^{u \varphi(q)} s ; g_{0}^{u \varphi(q)}\left(\alpha s+\beta\left(2 s-p^{t}\right)^{-1} s p^{t}\right)\right\}, \\
Q_{2}(s, u) & =\left(B_{2}^{(t)}(s)-p^{t}\right) \times g_{0}^{u \varphi(q)+\frac{\varphi(q)}{2}}, \\
& =\left\{p^{t}, 0, g_{0}^{\frac{\varphi(q)}{2}+u \varphi(q)}\left(s-p^{t}\right) ; g_{0}^{\frac{\varphi(q)}{2}+u \varphi(q)}\left(\alpha s+\beta\left(2 s-p^{t}\right)^{-1} s p^{t}-p^{t}\right)\right\}
\end{aligned}
$$

for $u \in\left[0, p^{t}-1\right]$ and $s \in S_{t}$, where the second equalities of $Q_{1}(s, u)$ and $Q_{2}(s, u)$ follow from $p^{t} g_{0}^{\varphi(q)} \equiv p^{t}\left(\bmod p^{m}\right)$ and $p^{t} g_{0}^{\frac{\varphi(q)}{2}} \equiv-p^{t}\left(\bmod p^{m}\right)$, respectively.

Let $H_{0}^{(t)}$ denote the multiplicative subgroup of $\mathbb{Z}_{p^{m}}^{\times}$generated by $g_{0}^{\varphi(q)}$ and denote the cosets of $H_{0}^{(t)}$ by $H_{a}^{(t)}=g_{0}^{a} H_{0}^{(t)}$.

First, consider the pure triples containing $\left\{0, p^{t}\right\}$ in $Q_{1}(s, u)$ and $Q_{2}(s, u)$. Let

$$
\begin{align*}
& U_{1}=\bigcup_{s \in S_{t}} \bigcup_{u=0}^{p^{t}-1}\left\{g_{0}^{u \varphi(q)} s\right\}=\bigcup_{s \in S_{t}} s H_{0}^{(t)}  \tag{2.31}\\
& U_{2}=\bigcup_{s \in S_{t}} \bigcup_{u=0}^{p^{t}-1}\left\{g_{0}^{\frac{\varphi(q)}{2}+u \varphi(q)}\left(s-p^{t}\right)\right\}=\bigcup_{s \in S_{t}}\left(s-p^{t}\right) H_{\frac{\varphi(q)}{2}}^{(t)} . \tag{2.32}
\end{align*}
$$

Furthermore, it follows from Proposition 2.3 .2 (iii) that,

$$
\begin{equation*}
U_{2}=\bigcup_{s \in S_{t}}\left(s-p^{t}\right) H_{\frac{\varphi(q)}{2}}^{(t)}=\bigcup_{s^{\prime} \in S_{t}} s^{\prime} H_{\frac{\varphi(q)}{2}}^{(t)} \tag{2.33}
\end{equation*}
$$

Then it follows from Proposition 2.3 .6 that $U_{1} \cup U_{2}=\mathbb{Z}_{p^{m}}^{\times}$. Therefore, the pure triple $\left\{0, p^{t}, x_{t}\right\}$ for every $x_{t} \in \mathbb{Z}_{p^{m}}^{\times}$is covered in $\bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{2}^{(t)}\left(s_{t}\right)\right)$.

Then, consider the mixed triples containing $\left\{0, p^{t}\right\}$ in $Q_{1}(s, u)$ and $Q_{2}(s, u)$. Let $\zeta(s)=s+\alpha^{-1} \beta\left(2 s-p^{t}\right)^{-1} s p^{t}$. Since $\zeta(s)$ is independent from $u$, we denote

$$
\begin{aligned}
& U_{1}(s)=\alpha \zeta(s) \bigcup_{u=0}^{p^{t}-1}\left\{g_{0}^{u \varphi(q)}\right\}=\alpha \zeta(s) H_{(t)}^{0} \\
& U_{2}(s)=\left(\alpha \zeta(s)-p^{t}\right) \bigcup_{u=0}^{p^{t}-1}\left\{g_{0}^{\frac{\varphi(q)}{2}+u \varphi(q)}\right\}=\left(\alpha \zeta(s)-p^{t}\right) H_{(t)}^{\frac{\varphi(q)}{2}}
\end{aligned}
$$

Furthermore, by Proposition 2.3.4, we have

$$
\begin{aligned}
& U_{1}=\bigcup_{s \in S_{t}} U_{1}(s)=\bigcup_{s^{\prime} \in S_{t}} \alpha s^{\prime} H_{(t)}^{0}, \\
& U_{2}=\bigcup_{s \in S_{t}} U_{2}(s)=\bigcup_{s^{\prime \prime} \in S_{t}} \alpha s^{\prime \prime} H_{(t)}^{\frac{\varphi(q)}{2}}
\end{aligned}
$$

It follows from Proposition 2.3.6 again that $\bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{2}^{(t)}\left(s_{t}\right)\right)$ covers the mixed triple $\left\{0, p^{t} ; y_{t}\right\}$ for every $y_{t} \in U_{1} \cup U_{2}=\alpha \mathbb{Z}_{p^{m}}^{\times}=\mathbb{Z}_{p^{m}}^{\times}$. In addition, since all the unions in this proof are between cosets, each pure or mixed triple we considered is covered exactly once.

Lemma 2.3.20 (Type $\left.^{\mathrm{III}}{ }^{\xi}\right)$. For $p \equiv 1(\bmod 12)$, let $B^{(t)}\left(s_{t}\right)=\left\{0,1, \xi_{m}+\right.$ $\left.s_{t} p^{t}, \overline{\xi_{m}}+\left(3-\sqrt{-3} s_{t} p^{t}\right)^{-1} s_{t} p^{t}\right\}$. Then, $\bigcup_{t=1}^{m-1} \bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B^{(t)}\left(s_{t}\right)\right)$ covers each pure triple of the form $\{0,1, x\}$ for $x \in\left(\bar{C}\left(\xi_{m}\right)+p \mathbb{Z}_{p^{m-1}}\right) \backslash \bar{C}\left(\xi_{m}\right)$ exactly once.

Proof. For a given $s$, denote

$$
x(s)=\xi_{m}+s p^{t} \quad \text { and } \quad y(s)=\overline{\xi_{m}}+\vartheta(s) p^{t}
$$

where $\vartheta(s)=\left(3-\sqrt{-3} s p^{t}\right)^{-1} s$. It can be verified that all pure triples containing $\{0,1\}$ covered by $\mathcal{O}_{A}\left(B^{(t)}(s)\right)$ are

$$
\left\{\{0,1, a\} \left\lvert\, a \in C(x(s)) \cup C(y(s)) \cup C\left(\frac{1-y(s)}{1-x(s)}\right) \cup C\left(\frac{x(s)-y(s)}{x(s)}\right)\right.\right\} .
$$

On the one hand, it is shown in Proposition 2.3.9 that

$$
\bigcup_{s \in R_{t}} C(x(s))=C\left(\xi_{m}\right)+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}
$$

On the other hand, we observe that

$$
\begin{equation*}
\left\{y(s), \frac{1-y(s)}{1-x(s)}, \frac{x(s)-y(s)}{x(s)}\right\}=\overline{\xi_{m}}+\left\{\vartheta(s) p^{t}, \xi_{m}^{2} \vartheta(s) p^{t}, \xi_{m}^{4} \vartheta(s) p^{t}\right\} \tag{2.34}
\end{equation*}
$$

For some $C(v)$ satisfying $C(v) \equiv C\left(\overline{\xi_{m}}\right)\left(\bmod p^{t}\right)$, each of the six elements of $C(v)$ are distinct by modulo $p^{t}$. Hence, for distinct $v_{1}, v_{2} \in \mathbb{Z}_{p^{m}}$, if $v_{1} \equiv v_{2} \equiv \overline{\xi_{m}}$
$\left(\bmod p^{t}\right)$, then $C\left(v_{1}\right) \equiv C\left(v_{2}\right) \equiv C\left(\overline{\xi_{m}}\right)\left(\bmod p^{t}\right)$ must hold, and $C\left(v_{1}\right)$ must be disjoint from $C\left(v_{2}\right)$. In particular, there is only one element $w$ in $C(v)$ such that $w \equiv \overline{\xi_{m}}\left(\bmod p^{t}\right)$. Now, we choose $y(s), \frac{1-y(s)}{1-x(s)}$ and $\frac{x(s)-y(s)}{x(s)}$ as representatives of their cross-ratio classes, and then show that

$$
\begin{equation*}
\bigcup_{s \in R_{t}}\left\{y(s), \frac{1-y(s)}{1-x(s)}, \frac{x(s)-y(s)}{x(s)}\right\}=\overline{\xi_{m}}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times} \tag{2.35}
\end{equation*}
$$

which implies that

$$
\bigcup_{s \in R_{t}}\left(C(y(s)) \cup C\left(\frac{1-y(s)}{1-x(s)}\right) \cup C\left(\frac{x(s)-y(s)}{x(s)}\right)\right)=C\left(\overline{\xi_{m}}\right)+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}
$$

2.35 can be easily shown by replacing the left-hand side with 2.34) and then applying Proposition 2.3.8. In addition, since $3\left|R_{t}\right|=\left|\mathbb{Z}_{p^{m-t}}^{\prime}\right|$, the number of appearances in $\mathcal{O}_{A}\left(B^{(t)}(s)\right)$ of every such pure triple is exactly one.

It remains to consider the pure and mixed triples containing $\left\{0, p^{t}\right\}$ for every $t \in[m-1]$. We complete this case by the following lemma without proof, since it is straightforward by the definition of an AsSQS.
Lemma 2.3.21 (Type V). Let $\left(\mathbb{Z}_{p^{m-1}} \times \mathbb{Z}_{2}, \mathcal{B}\right)$ be an AsSQS. For each $t \in$ $[m-1],\left\{p B\left(\bmod p^{m}\right) \mid B \in \mathcal{B}\right\}$ covers the pure triple $\left\{0, p^{t}, x_{t}\right\}$ for every $x_{t} \in p \mathbb{Z}_{p^{m-1}} \backslash\left\{0, p^{t}\right\}$ and the mixed triple $\left\{0, p^{t} ; y_{t}\right\}$ for every $y_{t} \in p \mathbb{Z}_{p^{m-1}}$.

We summarize the above lemmas of $\operatorname{AsSQS}^{A}\left(2 p^{m}\right)$ in Table 2.6. Table 2.7 and Table 2.8. In conclusion, if an $\operatorname{AsSQS}(2 p)$ exists, we can sequentially construct an $\operatorname{AsSQS}\left(2 p^{2}\right)$, an $\operatorname{AsSQS}\left(2 p^{3}\right), \ldots$, an $\operatorname{AsSQS}\left(2 p^{m}\right)$ for any positive integer $m$.

Table 2.6: Triples containing $\{0,1\}$ in an $\operatorname{AsSQS}^{A}\left(2 p^{m}\right)$ for $p \equiv 5(\bmod 12)$

| Type | Pure triples $\{0,1, x\}$, <br> for all $x$ in the following set | Mixed triples $\{0,1 ; y\}$, <br> for all $y$ in the following set | Lemmas |
| ---: | :--- | :--- | :--- |
| $\mathrm{I}^{\prime}$ |  | $\{\alpha, \beta\}$ | $\boxed{2.3 .13}$ |
| I |  | $\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi, 1-\chi\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\overline{2.3 .14}$ |
| II' | $\left\{-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\left\{0,1,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\underline{2.3 .15}$ |
| III | $\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ |  | $\underline{2.3 .17}$ |
| IV | $\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\{0,1\}$ | $\left(\{\alpha, \beta\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\{\alpha, \beta\}$ | $\overline{2.3 .18}$ |
| Union | $\mathbb{Z}_{p^{m}} \backslash\{0,1\}$ | $\mathbb{Z}_{p^{m}}$ |  |

### 2.3.3 Recursive construction $B$

In this subsection, we introduce a recursive construction for an $\operatorname{AsSQS}\left(2 p^{m}\right)$ based on an $\operatorname{AsSQS}^{B}(2 p)$, denoted by $\operatorname{AsSQS}^{B}\left(2 p^{m}\right)$. It is remarkable that any

Table 2.7: Triples containing $\{0,1\}$ in an $\operatorname{AsSQS}^{A}\left(2 p^{m}\right)$ for $p \equiv 1(\bmod 12)$

| Type | Pure triples $\{0,1, x\}$, for all $x$ in the following set | Mixed triples $\{0,1 ; y\}$, for all $y$ in the following set | Lemmas |
| :---: | :---: | :---: | :---: |
| $\mathrm{I}^{\prime}$ |  | $\{\alpha, \beta\}$ | 2.3.13 |
| I |  | $\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi, 1-\chi\right\}+p \mathbb{Z}_{p^{m-1}}$ | 2.3.14 |
| II' | $\left\{-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\left\{0,1,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | 2.3.15 |
| III | $\begin{aligned} \mathbb{Z}_{p} \backslash\left(\left\{0,1,-1,2,2^{-1}\right\}\right. & \left.\cup \bar{C}\left(\xi_{m}\right)\right) \\ & +p \mathbb{Z}_{p^{m-1}} \end{aligned}$ |  | 2.3.17 |
| III ${ }^{\xi}$ | $\bar{C}\left(\xi_{m}\right)+p \mathbb{Z}_{p^{m-1}}$ |  | 2.3.20 |
| IV | $\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\{0,1\}$ | $\left(\{\alpha, \beta\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\{\alpha, \beta\}$ | 2.3.18 |
| Union | $\mathbb{Z}_{p^{m}} \backslash\{0,1\}$ | $\mathbb{Z}_{p^{m}}$ |  |

Table 2.8: Triples containing $\left\{0, p^{t}\right\}$ in an $\operatorname{AsSQS}^{A}\left(2 p^{m}\right)$

| Type | Pure triples $\left\{0, p^{t}, x_{t}\right\}$, <br> for $t \in[m-1]$ | Mixed triples $\left\{0, p^{t} ; y_{t}\right\}$, <br> for $t \in[m-1]$ | Lemmas |
| ---: | :--- | :--- | :--- |
| IV | $x_{t} \in \mathbb{Z}_{p^{m}}^{\times}$ | $y_{t} \in \mathbb{Z}_{p^{m}}^{\times}$ | 2.3 .19 |
| V | $x_{t} \in p \mathbb{Z}_{p^{m-1}} \backslash\left\{0, p^{t}\right\}$ | $y_{t} \in p \mathbb{Z}_{p^{m-1}}$ | 2.3 .21 |
| Union | $x_{t} \in \mathbb{Z}_{p^{m}} \backslash\left\{0, p^{t}\right\}$ | $y_{t} \in \mathbb{Z}_{p^{m}}$ |  |

$\operatorname{AsSQS}^{B}(2 p)$ is 2-chromatic, accordingly, an $\operatorname{AsSQS}^{B}\left(2 p^{m}\right)$ is also 2-chromatic if all Type V base blocks in Construction 2.3 .22 come from an $\operatorname{AsSQS}^{B}\left(2 p^{m-1}\right)$.

Construction 2.3.22. Assume both an $\operatorname{AsSQS}^{B}(2 p)$ and an $\operatorname{AsSQS}^{B}\left(2 p^{m-1}\right)$ have been constructed, the base blocks of an $\operatorname{AsSQS}^{B}\left(2 p^{m}\right)$ can be obtained as follows:
(i) For the cases of prime $p \equiv 1$ or $5(\bmod 12)$, we have the base blocks as follows:

Type I': $\left\{0,1 ; \chi_{m}, 1-\chi_{m}\right\}$,
Type II': $\left\{0,1,-1+s p^{m-1} ; s p^{m-1}\right\}$ for $s \in\left[0, \frac{p-1}{2}\right]$,
Type IV: $\left\{0, p^{t}, s_{t} ; \chi_{m} s_{t}+\left(2 s_{t}-p^{t}\right)^{-1}\left(1-\chi_{m}\right) p^{t} s_{t}\right\}$ for $t \in[m-1]$ and $s_{t} \in S_{t}$, where $S_{t}$ is defined by (2.23),
Type V: $p B\left(\bmod p^{m}\right)$ for every base block $B$ of an $\operatorname{AsSQS}\left(2 p^{m-1}\right)$;
(ii) If $p \equiv 5(\bmod 12)$, we additionally have

Type II: $\left\{0,1, a_{i}+s p^{m-1} ; b_{i}+s p^{m-1}\right\}$ for $s \in[0, p-1]$ and $i \in\left[\frac{p-5}{6}\right] ;$
(iii) If $p \equiv 1(\bmod 12)$, we additionally have

Type $I I^{\xi}:\left\{0,1, \xi_{m} ; \overline{\xi_{m}}\right\}$ and $\left\{0,1, \xi_{m}+s_{t} p^{t} ; \overline{\xi_{m}}+\left(3-\sqrt{-3} s_{t} p^{t}\right)^{-1} s_{t} p^{t}\right\}$ for $t \in[m-1]$ and $s_{t} \in R_{t}$, where $R_{t}$ is defined by 2.25,

Table 2.9: Triples containing $\{0,1\}$ in an $\operatorname{AsSQS}^{B}\left(2 p^{m}\right)$ for $p \equiv 5(\bmod 12)$

| Type | Pure triples $\{0,1, x\}$, for all $x$ in the following set | Mixed triples $\{0,1 ; y\}$, for all $y$ in the following set | Lemmas |
| :---: | :---: | :---: | :---: |
| $\mathrm{I}^{\prime}$ |  | $\{\alpha, \beta\}$ | 2.3.13 |
| $\mathrm{II}^{\prime}$ | $\left\{-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\left\{0,1,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | 2.3.15 |
| II | $\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi, 1-\chi\right\}+p \mathbb{Z}_{p^{m-1}}$ | 2.3.23 (i) |
| IV | $\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\{0,1\}$ | $\left(\{\alpha, \beta\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\{\alpha, \beta\}$ | 2.3.18 |
| Union | $\mathbb{Z}_{p^{m}} \backslash\{0,1\}$ | $\mathbb{Z}_{p^{m}}$ |  |

Type II: $\left\{0,1, a_{i}+s p^{m-1} ; b_{i}+s p^{m-1}\right\}$ for $s \in[0, p-1]$ and $i \in\left[\frac{p-7}{6}\right]$.
Types I', II', IV, and V are exactly the same as those in Construction 2.3.10. hence we omit the proofs. For Type II base blocks, by a direct calculation which is analogous to the proof of Lemmas 2.3 .14 and 2.3 .17 , we state the following:

Lemma 2.3.23 (Type II). Let $B_{2}^{s}\left(a_{i}, b_{i}\right)=\left\{0,1, a_{i}+s p^{m-1} ; b_{i}+s p^{m-1}\right\}$.
(i) For $p \equiv 5(\bmod 12), \bigcup_{i=1}^{\frac{p-5}{6}} \bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{2}^{s}\left(a_{i}, b_{i}\right)\right)$ covers the pure triple $\{0,1, x\}$ for every

$$
x \in\left(\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}\right\}\right)+p \mathbb{Z}_{p^{m-1}}
$$

and the mixed triple $\{0,1 ; y\}$ for every

$$
y \in\left(\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi_{m}, 1-\chi_{m}\right\}\right)+p \mathbb{Z}_{p^{m-1}}
$$

(ii) For $p \equiv 1(\bmod 12), \bigcup_{i=1}^{\frac{p-7}{6}} \bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{2}^{s}\left(a_{i}, b_{i}\right)\right)$ covers the pure triple $\{0,1, x\}$ for every

$$
x \in\left(\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}, \xi_{m}, 1-\xi_{m}\right\}\right)+p \mathbb{Z}_{p^{m-1}}
$$

and the mixed triple $\{0,1 ; y\}$ for every

$$
y \in\left(\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi_{m}, 1-\chi_{m}, \overline{\xi_{m}}, 1-\overline{\xi_{m}}\right\}\right)+p \mathbb{Z}_{p^{m-1}}
$$

In summary, for $p \equiv 5(\bmod 12)$, Table 2.9 lists all types of base blocks whose orbits cover the triples of the form $\{0,1, x\}$ and $\{0,1 ; y\}$. The triples of the form $\left\{0, p^{t}, x_{t}\right\}$ and $\left\{0, p^{t} ; y_{t}\right\}$ are shown in Table 2.8 .

Moreover, for $p \equiv 1(\bmod 12), \mathcal{O}_{A}\left(\left\{0,1, \xi_{m} ; \overline{\xi_{m}}\right\}\right)$ covers pure triples $\left\{0,1, \xi_{m}\right\}$, $\left\{0,1,1-\xi_{m}\right\}$ and mixed triples $\left\{0,1 ; \overline{\xi_{m}}\right\},\left\{0,1 ; 1-\overline{\xi_{m}}\right\}$.

Lemma 2.3.24 (Type $\left.I I^{\xi}\right)$. Let $B^{(t)}\left(s_{t}\right)=\left\{0,1, \xi_{m}+s_{t} p^{t} ; \overline{\xi_{m}}+\left(3-\sqrt{-3} s_{t} p^{t}\right)^{-1} s_{t} p^{t}\right\}$
for each $s_{t} \in R_{t}$, where $R_{t}$ is defined as in 2.25). Then, $\bigcup_{t=1}^{m-1} \bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B^{(t)}\left(s_{t}\right)\right)$ covers the pure triple $\{0,1, x\}$ for every

$$
x \in\left\{\xi_{m}, 1-\xi_{m}\right\}+p \mathbb{Z}_{p^{m-1}} \backslash\left\{\xi_{m}, 1-\xi_{m}\right\}
$$

and the mixed triple $\{0,1 ; y\}$ for every

$$
y \in\left\{\overline{\xi_{m}}, 1-\overline{\xi_{m}}\right\}+p \mathbb{Z}_{p^{m-1}} \backslash\left\{\overline{\xi_{m}}, 1-\overline{\xi_{m}}\right\}
$$

exactly once.
Proof. Note that $\bigcup_{t=1}^{m-1} p^{t} \mathbb{Z}_{p^{m-t}}^{\times}=p \mathbb{Z}_{p^{m-1}} \backslash\{0\}$. Thus, it suffices to prove, for each $t \in[m-1], \bigcup_{s_{t} \in R_{t}} \mathcal{O}_{A}\left(B^{(t)}\left(s_{t}\right)\right)$ covers every triple in

$$
\begin{aligned}
& \left\{\{0,1 ; y\} \mid y \in\left\{\overline{\xi_{m}}, 1-\overline{\xi_{m}}\right\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}\right\} \cup \\
& \left\{\{0,1, x\} \mid x \in\left\{\xi_{m}, 1-\xi_{m}\right\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}\right\}
\end{aligned}
$$

exactly once. Denote $\vartheta\left(s_{t}\right)=\left(3-\sqrt{-3} s_{t} p^{t}\right)^{-1} s_{t}$ for $s_{t} \in R_{t}$. By tedious calculations we yield all the six quadruples containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B^{(t)}\left(s_{t}\right)\right)$, namely,

$$
\begin{align*}
& \left\{0,1, \xi_{m}+s_{t} p^{t} ; \overline{\xi_{m}}+\vartheta\left(s_{t}\right) p^{t}\right\}  \tag{2.36}\\
& \left\{0,\left(\xi_{m}+s_{t} p^{t}\right)^{-1}, 1 ; 1-\overline{\xi_{m}}-\xi_{m}^{4} \vartheta\left(s_{t}\right) p^{t}\right\}  \tag{2.37}\\
& \left\{1,1-\left(\xi_{m}+s_{t} p^{t}\right)^{-1}, 0 ; \overline{\xi_{m}}+\xi_{m}^{4} \vartheta\left(s_{t}\right) p^{t}\right\}  \tag{2.38}\\
& \left\{1,0,1-\xi_{m}-s_{t} p^{t} ; 1-\overline{\xi_{m}}-\vartheta\left(s_{t}\right) p^{t}\right\}  \tag{2.39}\\
& \left\{\left(1-\xi_{m}-s_{t} p^{t}\right)^{-1}, 0,1 ; \overline{\xi_{m}}+\xi_{m}^{2} \vartheta\left(s_{t}\right) p^{t}\right\}  \tag{2.40}\\
& \left\{1-\left(1-\xi_{m}-s_{t} p^{t}\right)^{-1}, 1,0 ; 1-\overline{\xi_{m}}-\xi_{m}^{2} \vartheta\left(s_{t}\right) p^{t}\right\} . \tag{2.41}
\end{align*}
$$

We first consider the mixed triples. For a fixed $s_{t}$, we can collect all the mixed triples contained in the quadruples (2.36)-2.41), given by

$$
\left\{\{0,1 ; y\} \mid y \in\left(\overline{\xi_{m}}+U\left(s_{t}\right)\right) \cup\left(1-\overline{\xi_{m}}-U\left(s_{t}\right)\right)\right\},
$$

where $U\left(s_{t}\right)=\left\{\vartheta\left(s_{t}\right), \xi_{m}^{2} \vartheta\left(s_{t}\right), \xi_{m}^{4} \vartheta\left(s_{t}\right)\right\}$. It follows from Proposition 2.3.8 that $\bigcup_{s_{t} \in R_{t}}\left(-U\left(s_{t}\right)\right)=\bigcup_{s_{t} \in R_{t}} U\left(s_{t}\right)=\mathbb{Z}_{p^{m-t}}^{\times}$. Therefore, for a certain $t \in[m-1]$, $\bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{s_{t}}^{(t)}\right)$ covers the mixed triple $\{0,1 ; y\}$ for every $y \in\left\{\overline{\xi_{m}}, 1-\overline{\xi_{m}}\right\}+$ $p^{t} \mathbb{Z}_{p^{m-t}}^{\times}$.

Next, all pure triples contained in the quadruples 2.36 - 2.41 are given by

$$
\left\{\{0,1, x\} \mid x \in C\left(\xi_{m}+s_{t} p^{t}\right)\right\}
$$

It follows from Proposition 2.3 .9 that

$$
\bigcup_{s_{t} \in R_{t}} C\left(\xi_{m}+s_{t} p^{t}\right)=\left\{\xi_{m}, 1-\xi_{m}\right\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}
$$

Therefore, for a certain $t, \bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B^{(t)}\left(s_{t}\right)\right)$ covers the pure triple $\{0,1, x\}$ for every $x \in\left\{\xi_{m}, 1-\xi_{m}\right\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}$.

In summary, for $p \equiv 1(\bmod 12)$, Table 2.10 lists all types of base blocks whose orbits cover the triples of the form $\{0,1, x\}$ and $\{0,1 ; y\}$.

Table 2.10: Triples containing $\{0,1\}$ in an $\operatorname{AsSQS}^{B}\left(2 p^{m}\right)$ for $p \equiv 1(\bmod 12)$

| Type | Pure triples $\{0,1, x\}$, for all $x$ in the following set | Mixed triples $\{0,1 ; y\}$, for all $y$ in the following set | Lemmas |
| :---: | :---: | :---: | :---: |
| $I^{\prime}$ |  | $\{\alpha, \beta\}$ | 2.3.13 |
| $\mathrm{II}^{\prime}$ | $\left\{-1,2,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\left\{0,1,2^{-1}\right\}+p \mathbb{Z}_{p^{m-1}}$ | 2.3.15 |
| II ${ }^{\xi}$ | $\left\{\xi_{m}, 1-\xi_{m}\right\}+p \mathbb{Z}_{p^{m-1}}$ | $\left\{\overline{\xi_{m}}, 1-\overline{\xi_{m}}\right\}+p \mathbb{Z}_{p^{m-1}}$ | 2.3.24 |
| II | $\mathbb{Z}_{p} \backslash\left\{0,1,-1,2,2^{-1}, \xi_{m}, 1-\xi_{m}\right\}$ | $\mathbb{Z}_{p} \backslash\left\{0,1,2^{-1}, \chi_{m}, 1-\chi_{m}, \overline{\xi_{m}}, 1-\overline{\xi_{m}}\right\}$ |  |
| IV | $\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\left\{\begin{array}{c} +p \mathbb{Z}_{p^{m-1}} \\ \{0,1\} \end{array}\right.$ | ${ }^{\left(\{\chi, 1-\chi\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash} \begin{aligned} & +p \mathbb{Z}_{p^{m-1}} \\ & \{\alpha, \beta\}\end{aligned}$ | $\frac{2.3 .23}{2.3 .18}(\mathrm{~b})$ |
| Union | $\mathbb{Z}_{p^{m}} \backslash\{0,1\}$ | $\mathbb{Z}_{p^{m}}$ |  |

### 2.4 A necessary condition for the existence of affine-invariant strictly cyclic Steiner quadruple systems

Recall that $n \equiv 1,5(\bmod 12)$ is necessary for the existence of an $\operatorname{sSQS}(2 n)$. A natural question arises: Is there any further requirement for an $\operatorname{AsSQS}(2 n)$ ? For example, it is known that an sSQS(98) exists (see 44 Example 7.7). Also, Example 2.2 .24 gives an "affine-invariant" $\operatorname{SQS}(98)$ over $\mathbb{F}_{49} \oplus \mathbb{F}_{2}$, which is not cyclic. Although we cannot construct an AsSQS(98) by our constructions, it would be interesting if it did exist. Now, we give a negative answer as follows:

Theorem 2.4.1. If there exists an $\operatorname{AsSQS}(2 n)$, then every prime factor $p$ of $n$ must satisfy $p \equiv 1,5(\bmod 12)$.

Proof. The necessary condition $n \equiv 1,5(\bmod 12)$ implies that all the prime factors of $n$ are congruent to $1,5,7,11$ modulo 12 . Thus it suffices to prove that $n$ does not have a prime factor $p$ with $p \equiv 7,11(\bmod 12)$.

Assume $n=p^{\alpha} q$, where $\alpha \geq 1$ and $q$ is coprime with $p$. We consider the quadruple containing $\{0, p, k p\}$ with $k \not \equiv 0,1\left(\bmod p^{\alpha-1} q\right)$. First, we suppose $B=\{0, p, k p, s\}$, where $s \not \equiv 0(\bmod p)$. Then, we can denote $s=a+b p$, for some $a \in \mathbb{Z}_{p}^{\times}$and $b \in \mathbb{Z}_{p^{\alpha-1} q}$. Let $\lambda=1+c p^{\alpha-1} q$ be an element in $\mathbb{Z}_{n}^{\times}$for some $c \in \mathbb{Z}_{p}^{\times}$. Then $\lambda p \equiv p(\bmod n)$ holds. Moreover, $\lambda s-s=$ $(\lambda-1) s=c p^{\alpha-1} q(a+b p) \equiv a c p^{\alpha-1} q \not \equiv 0(\bmod n)$ holds. Therefore, $B$ and $\lambda B$ are distinct and both of them contain the triple $\left\{0, p^{\alpha}, a p^{\alpha}\right\}$. Hence it suffices to consider the case when $B=\{0, p, k p, l p\}$, where $l \not \equiv 0,1, k\left(\bmod p^{\alpha-1} q\right)$. This is equivalent to saying $\{0,1, k, l\}$ is a base block of an $\operatorname{AsSQS}\left(2 p^{\alpha-1} q\right)$. By repeatedly applying this strategy, we can see that the existence of an $\operatorname{AsSQS}(2 n)$ requires that an $\operatorname{AsSQS}(2 p)$ exists for every prime divisor $p$ of $n$.

As an open problem and for future work, we are also interested in the existence of an affine-invariant $\operatorname{SQS}\left(2^{\alpha_{0}} p_{1}^{\alpha_{1}} p_{2}^{\alpha_{2}} \cdots p_{r}^{\alpha_{r}}\right)$ which is not necessarily strictly cyclic, where $p_{i} \equiv 1,5(\bmod 12)$ is prime for every $i \in[r]$. It is also a
challenge to consider packing designs and covering designs in the same manner for applications.

### 2.5 Affine-invariant two-fold quadruple systems over $\mathbb{Z}_{p}$

This section is devoted to providing a direct construction of an affine-invariant TQS (two-fold quadruple system) via the graphs CG $\left(\Omega_{p}\right)$. Roughly speaking, by removing a 1 -factor from $\mathrm{CG}\left(\Omega_{p}\right)$, the resulting graph leads to the base blocks of an affine-invariant TQS $(p)$.
Construction 2.5.1. For prime $p \equiv 5(\bmod 12)$, suppose $\mathrm{CG}\left(\Omega_{p}\right)$ has a 1 factor, say $F$. Let $\ell=\frac{p-5}{6}$ and let $a_{1}, a_{2}, \ldots, a_{\ell}$ be elements in $\Omega_{q}$ such that

$$
\left\{\left\{C\left(a_{1}\right), C\left(\overline{a_{1}}\right)\right\},\left\{C\left(a_{2}\right), C\left(\overline{a_{2}}\right)\right\}, \ldots,\left\{C\left(a_{\ell}\right), C\left(\overline{a_{\ell}}\right)\right\}\right\}=E\left(\mathrm{CG}\left(\Omega_{q}\right)\right) \backslash E(F)
$$

where $E\left(\mathrm{CG}\left(\Omega_{q}\right)\right)$ and $E(F)$ denote the edge set of $\mathrm{CG}\left(\Omega_{q}\right)$ and $F$, respectively. Then

$$
\mathcal{B}=\left\{B_{a_{i}}=\left\{0,1, a_{i}, 1-a_{i}\right\} \mid i \in[\ell]\right\} \cup\left\{B_{-1}\right\} .
$$

is the set of base blocks of an affine-invariant $\operatorname{TQS}(p)$.
Theorem 2.5.2. An affine-invariant $\operatorname{TQS}(p)$ exists if the graph $\operatorname{CG}\left(\Omega_{p}\right)$ has a 1-factor.

Proof. As shown in the proof of Lemma 2.2.4 for Type III base blocks of $\operatorname{AsSQS}^{A}(2 p), \mathcal{O}_{A}\left(B_{a_{i}}\right)$ covers all triples of the form $\{0,1, x\}$ for $x \in C\left(a_{i}\right) \cup C\left(\overline{a_{i}}\right)$ if $C\left(a_{i}\right) \neq C\left(\overline{a_{i}}\right)$. It is shown in Proposition 2.1.10 (iii) that all the vertices of $\mathrm{CG}\left(\Omega_{q}\right)$ are of degree 3 except $C(3)$ and $C(\mu)$ for $q \equiv 5(\bmod 12)$. Hence, the multiset union of all the non-loop edges in $E\left(\operatorname{CG}\left(\Omega_{q}\right)\right) \backslash E(F)$ covers all the vertices of $\mathrm{CG}\left(\Omega_{q}\right)$ twice except $C(3), C(\chi)$, and $C(\mu)$, where each of $C(3)$ and $C(\chi)$ occurs once, and $C(\mu)$ does not appear.

Moreover, by Proposition 2.1.10 (i), $C\left(a_{i}\right)=C\left(\overline{a_{i}}\right)$ if and only if $C\left(a_{i}\right)=$ $C(\chi)$ or $C(\mu)$, where $\chi=\frac{1+\sqrt{ }-1}{2}$ satisfies $\chi=\chi^{\sigma_{C}}$ and $\mu=\frac{3+\sqrt{5}}{2}$ satisfies $\mu^{\sigma_{B}}=\mu^{\sigma_{C}}$ (cf. Table 2.1). In these cases, we can explicitly derive all the blocks containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{\chi}\right)$ and $\mathcal{O}_{A}\left(B_{\mu}\right)$, namely,

$$
\begin{aligned}
& \{0,1, \chi, 1-\chi\},\left\{0, \frac{1}{\chi}, 1, \frac{\chi}{\chi-1}\right\},\left\{1, \frac{\chi-1}{\chi}, 0, \frac{1}{1-\chi}\right\}, \text { and } \\
& \{0,1, \mu, 1-\mu\},\left\{0, \frac{1}{\mu}, 1, \frac{1}{1-\mu}\right\},\left\{1, \frac{\mu-1}{\mu}, 0, \frac{\mu}{\mu-1}\right\}, \\
& \left\{0, \frac{1}{1-\mu}, 1-\mu, 1\right\},\left\{1, \frac{\mu}{\mu-1}, \mu, 0\right\},\left\{\frac{1}{\mu}, \frac{\mu-1}{\mu}, 1,0\right\},
\end{aligned}
$$

in which every element of $C(\chi)$ occurs once and every element of $C(\mu)$ occurs twice.

Last, all the blocks containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{-1}\right)$ are

$$
\{0,1,-1,2\},\{0,-1,1,-2\},\left\{0, \frac{1}{2},-\frac{1}{2}, 1\right\},\left\{1,2^{-1}, \frac{3}{2}, 0\right\},\{2,1,3,0\},\left\{\frac{2}{3}, \frac{1}{3}, 1,0\right\}
$$

where each element of $C(2)$ occurs twice and each element of $C(3)$ occurs once.
Summing up the elements being covered in the above three cases, we observe that every element of $\Omega_{q} \cup C(2)$, that is $\mathbb{F}_{q} \backslash\{0,1\}$, appears twice. Therefore, $\bigcup_{B \in \mathcal{B}} \mathcal{O}_{A}(B)$ covers all triples of the form $\{0,1, x\}$ twice for every $x \in \mathbb{F}_{q} \backslash$ $\{0,1\}$.

Example 2.5.3. Let $p=29$. Then $B_{-1}, B_{14}, B_{4}, B_{25}, B_{9}$ are the base blocks of an affine-invariant TQS $(p)$, where the last four blocks respectively correspond to the edges $\{C(3), C(4)\},\{C(4), C(9)\},\{C(5)\}$ (a self-loop), $\{C(9)\}$ (a self-loop) of $\mathrm{CG}\left(\Omega_{29}\right)$ illustrated in Figure 2.5 (cf. Figure 2.2 ).

### 2.6 Affine-invariant two-fold quadruple systems over $\mathbb{Z}_{p^{m}}$

In this section, we construct an affine-invariant $\operatorname{TQS}\left(p^{m}\right)$ via the affine-invariant TQS $(p)$ obtained from Construction 2.5.1

Let $\chi_{t}$ denote a root of $2 \chi_{t}^{2}-2 \chi_{t}+1=0$ over $\mathbb{Z}_{p^{t}}$ for $t \in[1, m]$. Let $\mu_{t}$ denote a root of $\mu_{t}^{2}-3 \mu_{t}+1=0$ over $\mathbb{Z}_{p^{t}}$ for $t \in[1, m]$. Let $B_{s}(a)=$ $\left\{0,1, a+s p^{m-1}, 1-\left(a+s p^{m-1}\right)\right\}$.
Construction 2.6.1. Suppose $p \equiv 5(\bmod 12)$ is prime. We use the same notation with Construction 2.5.1. Assume that both an affine-invariant TQS $(p)$ and an affine-invariant TQS $\left(p^{m-1}\right)$ have been constructed, then the base blocks of an affine-invariant $\operatorname{TQS}\left(p^{m}\right)$ can be obtained as follows:

Type I: $B_{s}\left(a_{i}\right)$ for $i \in\left[\frac{p-5}{6}\right]$ and $s \in[0, p-1]$;
Type II: $B_{s}(-1)$ for $s \in[0, p-1]$;
Type III: $B_{s}\left(\chi_{m}\right)$ for $s \in\left[0, \frac{p-1}{2}\right]$;
Type IV: $B_{s}\left(\mu_{m}\right)$ for $s \in[0, p-1]$, if $p \equiv 29,41(\bmod 60)$;
Type $\mathrm{V}:\left\{0, p^{t}, s_{t}, s_{t}+p^{t}\right\}$, for $t \in[m-1]$ and $s_{t} \in S_{t}$, where $S_{t}$ is defined by (2.23);

Type VI: $p B\left(\bmod p^{m}\right)$, for all base blocks $B$ of the affine-invariant TQS $\left(p^{m-1}\right)$.
Lemma 2.6.2. For a fixed $a \in \Omega_{p} \backslash\left(C\left(\chi_{1}\right) \cup C\left(\mu_{1}\right)\right), \bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{s}(a)\right)$ covers $\{0,1, x\}$ exactly once for every $x \in \bar{C}(a)+p \mathbb{Z}_{p^{m-1}}$.
Proof. It suffices to show $\bigcup_{s=0}^{p-1} \bar{C}\left(a+s p^{m-1}\right)=\bar{C}(a)+p \mathbb{Z}_{p^{m-1}}$. First, $C(a+$ $\left.s p^{m-1}\right)$ and $C\left(\overline{a+s p^{m-1}}\right)$ are clearly disjoint. By Lemma 2.3.16. $\bar{C}\left(a+s_{1} p^{m-1}\right)$ and $\bar{C}\left(a+s_{2} p^{m-1}\right)$ are disjoint if $s_{1} \neq s_{2}$. Therefore, we have

$$
\bigcup_{s=0}^{p-1} \bar{C}\left(a+s p^{m-1}\right) \subseteq \bar{C}(a)+p \mathbb{Z}_{p^{m-1}}
$$

Furthermore, both the left-hand side and the right-hand side have cardinalities $12 p^{m-1}$, which completes the proof.

In the same manner, we can easily obtain the following lemmas.
Lemma 2.6.3. $\bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{s}(-1)\right)$ covers $\{0,1, x\}$ once for every $x \in C(3)+$ $p \mathbb{Z}_{p^{m-1}}$ and twice for every $x \in C(2)+p \mathbb{Z}_{p^{m-1}}$.

Proof. It is easy to check that $\mathcal{O}_{A}\left(B_{s}(-1)\right)$ covers $\{0,1, x\}$ for $x \in C(-1+$ $\left.s p^{m-1}\right) \cup C\left(3+s p^{m-1}\right)$. Then we can complete the proof by noting that $C(-1+$ $\left.s p^{m-1}\right)=\left\{-1 \pm s p^{m-1}, 2 \pm s p^{m-1}, 2^{-1} \pm 2^{-1} s p^{m-1}\right\}$ and $C\left(3+s p^{m-1}\right)=$ $\left\{3+s p^{m-1},-2-s p^{m-1},-\frac{1}{2}+\frac{1}{4} s p^{m-1}, \frac{3}{2}-\frac{1}{4} s p^{m-1}, \frac{2}{3}+\frac{1}{9} s p^{m-1}, \frac{1}{3}-\frac{1}{9} s p^{m-1}\right\}$.

Lemma 2.6.4. $\bigcup_{s=0}^{\frac{p-1}{2}} \mathcal{O}_{A}\left(B_{s}\left(\chi_{m}\right)\right)$ covers $\{0,1, x\}$ once for every $x \in C\left(\chi_{1}\right)+$ $p \mathbb{Z}_{p^{m-1}}$.

Proof. Note that $\overline{\chi_{m}+s p^{m-1}}=\chi_{m}-s p^{m-1}$. Hence

$$
\bigcup_{s=1}^{\frac{p-1}{2}}\left\{C\left(\chi_{m}+s p^{m-1}\right), C\left(\overline{\chi_{m}+s p^{m-1}}\right)\right\}=\bigcup_{s=1}^{p-1}\left\{C\left(\chi_{m}+s p^{m-1}\right)\right\}
$$

Therefore, $\bigcup_{s=1}^{\frac{p-1}{2}} \mathcal{O}_{A}\left(B_{s}\left(\chi_{m}\right)\right)$ covers $\{0,1, x\}$ once for every $x \in C\left(\chi_{m}\right)+$ $p\left(\mathbb{Z}_{p^{m-1}} \backslash\{0\}\right)$. In addition, it is known that $\mathcal{O}_{A}\left(B_{0}\left(\chi_{m}\right)\right)$ covers $\{0,1, x\}$ once for every $x \in C\left(\chi_{m}\right)$. Thus, $\bigcup_{s=0}^{\frac{p-1}{2}} \mathcal{O}_{A}\left(B_{s}\left(\chi_{m}\right)\right)$ covers $\{0,1, x\}$ once for every $x \in C\left(\chi_{m}\right)+p \mathbb{Z}_{p^{m-1}}=C\left(\chi_{1}\right)+p \mathbb{Z}_{p^{m-1}}$.
Lemma 2.6.5. $\bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{s}\left(\mu_{m}\right)\right)$ covers $\{0,1, x\}$ twice for every $x \in C\left(\mu_{1}\right)+$ $p \mathbb{Z}_{p^{m-1}}$.

Proof. Note that $C\left(\overline{\mu_{m}+s p^{m-1}}\right)=C\left(\mu_{m}-\mu_{m}^{-1} s p^{m-1}\right)$. Moreover, $\mathcal{O}_{A}\left(B_{0}\left(\mu_{m}\right)\right)$ covers $\{0,1, x\}$ twice for each $x \in C\left(\mu_{m}\right)$. Therefore, $\bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{s}\left(\mu_{m}\right)\right)$ covers $\{0,1, x\}$ twice for every $x \in C\left(\mu_{m}\right)+p \mathbb{Z}_{p^{m-1}}=C\left(\mu_{1}\right)+p \mathbb{Z}_{p^{m-1}}$.

Lemma 2.6.6 (Type I, II, III, IV). Let $\mathcal{O}_{1}=\bigcup_{i=1}^{\frac{p-5}{6}} \bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{s}\left(a_{i}\right)\right), \mathcal{O}_{2}=$ $\bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{s}(-1)\right)$, $\mathcal{O}_{3}=\bigcup_{s=0}^{\frac{p-1}{2}} \mathcal{O}_{A}\left(B_{s}\left(\chi_{m}\right)\right)$, and $\mathcal{O}_{4}=\bigcup_{s=0}^{p-1} \mathcal{O}_{A}\left(B_{s}\left(\mu_{m}\right)\right)$. Then $\bigcup_{k=1}^{4} \mathcal{O}_{k}$ covers $\{0,1, x\}$ twice for every $x \in\left(\mathbb{Z}_{p} \backslash\{0,1\}\right)+p \mathbb{Z}_{p^{m-1}}$.

Proof. It follows from Construction 2.5 .1 and Lemma 2.6 .2 that $\mathcal{O}_{1}$ covers $\{0,1, x\}$ exactly twice for every $x \in\left(\mathbb{Z}_{p} \backslash\left(C(0) \cup C(2) \cup C(3) \cup C\left(\chi_{1}\right) \cup\right.\right.$ $\left.C\left(\mu_{1}\right)\right)+p \mathbb{Z}_{p^{m-1}}$ and exactly once for $x \in C(3)+p \mathbb{Z}_{p^{m-1}}$. It can be immediately concluded, by combining Lemmas 2.6.3 2.6.4 and 2.6.5, that $\bigcup_{k=1}^{4} \mathcal{O}_{k}$ covers $\{0,1, x\}$ twice for every $x \in\left(\mathbb{Z}_{p} \backslash\{0,1\}\right)+p \mathbb{Z}_{p^{m-1}}$.

Lemma 2.6.7 (Type V). Let $B_{s_{t}}^{(t)}=\left\{0, p^{t}, s_{t}, s_{t}+p^{t}\right\}$. Then $\bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{s_{t}}^{(t)}\right)$ covers $\{0,1, x\}$ exactly twice for every $\left.x \in\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right\}\right) \backslash\{0,1\}$.

Proof. There are four distinct blocks containing $\{0,1\}$ in $\mathcal{O}_{A}\left(B_{s_{t}}^{(t)}\right)$.

$$
\begin{align*}
B^{\prime} & =B_{s_{t}}^{(t)} \times s_{t}^{-1}=\left\{0, s_{t}^{-1} p^{t}, 1,1+s_{t}^{-1} p^{t}\right\}  \tag{2.42}\\
\tilde{B}^{\prime} & =B^{\prime} \times\left(1+s_{t}^{-1} p^{t}\right)^{-1}=\left\{0, p^{t}\left(s_{t}+p^{t}\right)^{-1},\left(1+s_{t}^{-1} p^{t}\right)^{-1}, 1\right\}  \tag{2.43}\\
B^{\prime \prime} & =\left(B_{s_{t}}^{(t)}-p^{t}\right) \times\left(s_{t}-p^{t}\right)^{-1}=\left\{p^{t}\left(-s_{t}+p^{t}\right)^{-1}, 0,1,\left(1-s_{t}^{-1} p^{t}\right)^{-1}\right\}  \tag{2.44}\\
\tilde{B}^{\prime \prime} & =B^{\prime \prime} \times\left(1-s_{t}^{-1} p^{t}\right)=\left\{-s_{t}^{-1} p^{t}, 0,1-s_{t}^{-1} p^{t}, 1\right\} \tag{2.45}
\end{align*}
$$

where $\tilde{B}^{\prime \prime}=1-B^{\prime}, \tilde{B}^{\prime}=1-\tilde{B}^{\prime}$, and $B^{\prime \prime}=1-B^{\prime \prime}$. Moreover, we have $\left(1 \pm s_{t}^{-1} p^{t}\right)^{-1}-1=-p^{t}\left( \pm s_{t}+p^{t}\right)^{-1}$. Recall Proposition 2.3.2 (iii) and (iii) that $S_{t}+p^{t} \equiv S_{t}\left(\bmod p^{m-t}\right)$ and $S_{t} \cup\left(-S_{t}\right)=\mathbb{Z}_{p^{m-t}}^{\times}$. Hence,

$$
\begin{aligned}
& \left\{s_{t}^{-1} \mid s_{t} \in S_{t}\right\} \cup\left\{-s_{t}^{-1} \mid s_{t} \in S_{t}\right\}=\mathbb{Z}_{p^{m-t}}^{\times} \text {and } \\
& \left\{\left(s_{t}+p^{t}\right)^{-1} \mid s_{t} \in S_{t}\right\} \cup\left\{\left(-s_{t}+p^{t}\right)^{-1} \mid s_{t} \in S_{t}\right\}=\mathbb{Z}_{p^{m-t}}^{\times}
\end{aligned}
$$

Therefore, for any given $t \in[m-1]$, the union of all the triples containing $\{0,1\}$ in 2.42-2.45 extended over $s_{t} \in S_{t}$ covers $\left\{\{0,1, x\} \mid x \in\{0,1\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}\right\}$ twice. Furthermore, by $\left.\bigcup_{t=1}^{m-1}\left(\{0,1\}+p^{t} \mathbb{Z}_{p^{m-t}}^{\times}\right)=\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right\}\right) \backslash\{0,1\}$, the proof is completed.

Lemma 2.6.8 (Type V). For any fixed $t \in[m-1], \bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{s_{t}}^{(t)}\right)$ covers $\left\{0, p^{t}, x\right\}$ exactly twice for every $x \in \mathbb{Z}_{p^{m}}^{\times}$.

Proof. The idea is the same as the proof of Lemma 2.3 .19 for Type IV base blocks of $\operatorname{AsSQS}^{A}\left(2 p^{m}\right)$. Let $g_{0}$ be a generator of $\mathbb{Z}_{p^{m}}^{\times}$and simply denote $q=$ $p^{m-t}$. For a given $s \in S_{t}$, we can derive the blocks containing $\left\{0, p^{t}\right\}$ in $\mathcal{O}_{A}\left(B_{s}^{(t)}\right)$ as follows: Let
$Q_{1}(s, u)=B_{2}^{(t)}(s) \times g_{0}^{u \varphi(q)}=\left\{0, p^{t}, g_{0}^{u \varphi(q)} s, p^{t}+g_{0}^{u \varphi(q)} s\right\}$ and
$Q_{2}(s, u)=\left(B_{2}^{(t)}(s)-p^{t}\right) \times g_{0}^{u \varphi(q)+\frac{\varphi(q)}{2}}=\left\{p^{t}, 0, p^{t}+g_{0}^{\frac{\varphi(q)}{2}+u \varphi(q)} s, g_{0}^{\frac{\varphi(q)}{2}+u \varphi(q)} s\right\}$
for $u \in\left[0, p^{t}-1\right]$, which follow from $p^{t} g_{0}^{\varphi(q)} \equiv p^{t}\left(\bmod p^{m}\right)$ and $p^{t} g_{0}^{\frac{\varphi(q)}{2}} \equiv-p^{t}$ $\left(\bmod p^{m}\right)$, respectively. We have

$$
\bigcup_{s \in S_{t}} \bigcup_{u=0}^{p^{t}-1}\left\{g_{0}^{u \varphi(q)} s, g_{0}^{\frac{\varphi(q)}{2}+u \varphi(q)} s\right\}=\bigcup_{s \in S_{t}} s\left(H_{0}^{(t)} \cup H_{\frac{\varphi(q)}{2}}^{(t)}\right)=\mathbb{Z}_{p^{m}}^{\times}
$$

where the last equality follows from Proposition 2.3.6. Hence, for every $x \in \mathbb{Z}_{p^{m}}^{\times}$, the triple $\left\{0, p^{t}, x\right\}$ is covered twice in $\bigcup_{s_{t} \in S_{t}} \mathcal{O}_{A}\left(B_{s_{t}}^{(t)}\right)$.

Let $p \mathcal{B}_{p}=\left\{p B\left(\bmod p^{m}\right) \mid B \in \mathcal{B}_{p}\right\}$, where $\mathcal{B}_{p}$ denotes the set of all base blocks of the affine-invariant $\operatorname{TQS}\left(p^{m-1}\right)$ obtained from Construction 2.5.1

Lemma 2.6.9 (Type VI). For each $t \in[m-1]$, $\bigcup_{B^{*} \in p \mathcal{B}_{p}} \mathcal{O}_{A}\left(B^{*}\right)$ covers $\left\{0, p^{t}, x_{t}\right\}$ for every $x_{t} \in p \mathbb{Z}_{p^{m-1}} \backslash\left\{0, p^{t}\right\}$ exactly twice.
Proof. This is obvious from Construction 2.5.1.
We can sum up Lemmas 2.6.6, 2.6.8, 2.6.7 and 2.6.9 as follows:
(i) (Lemma 2.6.6 The affine orbits of Type I, II, III, and IV blocks cover $\{0,1, x\}$ twice for every $x \in \mathbb{Z}_{p} \backslash\{0,1\}+p \mathbb{Z}_{p^{m-1}}$.
(ii) (Lemma 2.6.7) The affine orbits of Type V blocks cover $\{0,1, x\}$ twice for every $x \in\left(\{0,1\}+p \mathbb{Z}_{p^{m-1}}\right) \backslash\{0,1\}$.
(iii) (Lemma 2.6.8 The affine orbits of Type V blocks cover $\left\{0, p^{t}, x\right\}$ twice for each $t \in[m-1]$ and every $x \in \mathbb{Z}_{p^{m}}^{\times}$.
(iv) (Lemma 2.6.9) The affine orbits of Type V blocks cover $\{0, p, x\}$ twice for each $t \in[m-1]$ and every $x_{t} \in p \mathbb{Z}_{p^{m-1}} \backslash\left\{0, p^{t}\right\}$.

In summary, we have the following criterion for the existence of affineinvariant TQS.

Theorem 2.6.10. For $p \equiv 5(\bmod 12)$, if the graph $\mathrm{CG}\left(\Omega_{p}\right)$ has no bridge besides its pendant edge, then an affine-invariant $\operatorname{TQS}\left(p^{m}\right)$ exists for any positive integer $m$.

### 2.7 Applications

In addition to the applications to OOCs, we briefly introduce the merits of the affine-invariant property for other applications. In particular, the constructions of 3 -designs are usually more complicated than 2 -designs. It is necessary to consider the practical significance for applications.

### 2.7.1 Searching blocks

Let $(V, \mathcal{B})$ be a $t$-design. For a given $s$-subset with $s \leq t$ (for example, a pair, a triple, etc.) of the point set $V$, say $T$, it is usually required to find the blocks containing a certain $T$ in the applications to group testing 39, filing schemes [6, 126], etc. The affine-invariant property works effectively for these kind of problems.

We recall the AsSQS(25) in Example 2.3.11 to illustrate the main idea for searching blocks. First, we index every base block as shown in Table 2.11. Then, we create a "retrieval table" for the triples containing $\{0,1\}$ as shown in Table 2.12 For example, when the quadruple containing $\{1,3 ; 5\}$, say $Q$, is requested, we follow these steps:

Step 1: $\quad$ Since $(\{1,3 ; 5\}-1) \times 2^{-1}=\{0,1 ; 2\}$, then $\mathcal{O}_{A}(\{1,3 ; 5\})=\mathcal{O}_{A}(\{0,1 ; 2\})$.

Table 2.11: Index of base blocks

| No. | Base blocks | No. | Base blocks |
| :--- | :--- | :--- | :--- |
| 1 | $\{0,1 ; 4,22\}$ | 5 | $\{0,5,1 ; 9\}$ |
| 2 | $\{0,1,24 ; 0\}$ | 6 | $\{0,5,2 ; 13\}$ |
| 3 | $\{0,1,4 ; 5\}$ | 7 | $\{0,5 ; 10,20\}$ |
| 4 | $\{0,1,9 ; 10\}$ | 8 | $\{0,5,20 ; 0\}$ |

Table 2.12: Retrieval table of cyclic orbits in each affine orbit

| Triples | No. of base blocks | The other element |
| :--- | :--- | :--- |
| $\{0,1 ; 0\}$ | 2 | 24 |
| $\{0,1 ; 1\}$ | 2 | 2 |
| $\{0,1 ; 2\}$ | 5 | 6 |
| $\{0,1 ; 3\}$ | 4 | 23 |
| $\vdots$ | $\vdots$ | $\vdots$ |

Step 2: Find $\{0,1 ; 2\}$ in Table 2.12 and observe $\{0,1 ; 2\}$ is covered by the affine orbit of No. 5 base block, i.e., $\mathcal{O}_{A}(\{0,5,1 ; 9\})$. Then, the desired quadruple for $\{0,1 ; 2\}$ is $\{0,1,6 ; 2\}$.

Step 3: Take $Q=\{0,1,6 ; 2\} \times 2+1=\{1,3,13 ; 5\}$.
In summary, we first find the desired "affine orbits", then the "cyclic orbits", and finally the "blocks". This procedure can be regarded as a generalization of the "retrieval algorithm" in filing schemes by using cyclic 2-designs (difference families) (see [6]). Clearly, by making use of the structure of automorphism groups, it is much more efficient than searching among all blocks.

### 2.7.2 Generating blocks

It is usually desired to generate a certain part of blocks for the applications of authentication codes [88, 110] as fast as possible using less storage. The structure of "affine orbits - cyclic orbits - blocks" also helps us to avoid unnecessary computation. Moreover, the storage requirement of affine base blocks of an $\operatorname{AsSQS}(v)$ is approximately reduced by up to $O(v)$ times from that of cyclic base blocks (see Tables $2.2,2.3,2.4$, and 2.5 .

On the other hand, practical applications of authentication codes often ask for an extremal large-scale design. The procedure of generating base blocks of an $\operatorname{AsSQS}(2 p)$ of Construction 2.2 .6 relies on a 1 -factor of the graph $\operatorname{CG}\left(\Omega_{p}\right)$. Clearly, if a 1-factor is known, it needs at most $O(p)$ time to generate all base blocks. Note that $\mathrm{CG}\left(\Omega_{p}\right)$ is an "almost" 3-regular graph of order $O(p)$ by Proposition 2.1.10 (iii). By using an algorithm for the maximum matching problem by Micali and Vazirani (see [91, 114), it can be completed in $O\left(p^{3 / 2}\right)$ time
to find a 1-factor of $\operatorname{CG}\left(\Omega_{p}\right)$. Actually, if we assume that $\mathrm{CG}\left(\Omega_{p}\right)$ has no bridge except its pendant edges, by using Diks and Stańczyk's algorithm 38 for a 2 -connected 3 -regular graph, a 1 -factor can be found in $O\left(p \log ^{2} p\right)$ time. In summary, even if all the blocks (the whole incident matrix) are required, it can be done in $O\left(p^{3 / 2}\right)$ time.

## Chapter 3

## Grid-block difference families

This chapter is devoted to the existence and construction of grid-block difference families, which can be regarded as generalizations of difference families and cyclic grid-block designs.

In order to show the existence, we first introduce an intermediate conclusion for estimating the asymptotic existence of an element satisfying certain cyclotomic conditions in a finite field.

### 3.1 An intermediate consequence derived from Weil's Theorem on multiplicative character sums

For many direct constructions of designs, the essential problem often comes down to choosing a proper subset to form an SDR (system of distinct representatives) of a certain set system. For "DF-like" structures over $\mathbb{F}_{q}$, the "certain set system" is usually a collection of cyclotomic cosets. In this case, Buratti and Pasotti's Theorem 1.4.8 provides a general solution when the desired "proper subset" forms a system of linear functions with respect to some $x \in \mathbb{F}_{q}$.

For instance, for showing the existence of a $(q, 6,1)$-DF over $\mathbb{F}_{q}$, Wilson proposed a sufficient condition as follows:

Theorem 3.1.1 (Wilson [120] Theorem 11). Let $q \equiv 1(\bmod 30)$ be a prime power and let $\omega$ be a primitive cube root of unity in $\mathbb{F}_{q}$. If there exists an element $x \in \mathbb{F}_{q}$ such that $\left\{1, x, \frac{x-1}{\omega-1}, \frac{x-\omega}{\omega-1}, \frac{x-\omega^{2}}{\omega-1}\right\}$ forms a system of representatives for $\mathcal{C}^{(5)}$, then there exists a $(q, 6,1)$-DF over $\mathbb{F}_{q}$.

In order to meet the above criteria, one can suppose

$$
x \in C_{1}^{(5)}, \frac{x-1}{\omega-1} \in C_{2}^{(5)}, \frac{x-\omega}{\omega-1} \in C_{3}^{(5)}, \text { and } \frac{x-\omega^{2}}{\omega-1} \in C_{4}^{(5)}
$$

and use Theorem 1.4 .8 to obtain a bound on $q$. Actually, Chen and Zhu 28 considered a more relaxed assumption, say

$$
x \in C_{i}^{(5)}, \frac{x-1}{\omega-1} \in C_{2 i}^{(5)}, \frac{x-\omega}{\omega-1} \in C_{3 i}^{(5)}, \text { and } \frac{x-\omega^{2}}{\omega-1} \in C_{4 i}^{(5)} \text { for some } i \in \mathbb{Z}_{5} \backslash\{0\}
$$

where the subscripts of $C_{i j}^{(5)}$ are reduced modulo 5 . Consequently, a better bound on $q$ was obtained.

The same trick has also been used by Chen and Zhu [30 to improve the existence bound for a $(q, 7,1)$-DF in which every base block is of the form $\left\{0,1, x, x^{2}, x^{3}, x^{4}, x^{5}\right\}$.

In general, if a complete $S D R$ of $\mathcal{C}^{(e)}$ is desired instead of an "SDR", any $m$ in the group of units $\mathbb{Z}_{e}^{\times}$can be used as a multiplier to the subscripts. Then, a better existence bound can be obtained from the relaxed criteria.

Note that, for instance, if $\left(1,2, x_{2}, x_{3}, x_{4}\right)$ is desired to form a complete system of representatives of $\mathcal{C}^{(5)}$, this idea does not work well anymore (see, for instance, Chen, Wei, and Zhu [27] on ( $q, 7,1$ )-DF). But in this case, Buratti and Pasotti's Theorem 1.4.8 can also give an answer.

In what follows, let $\mu(\cdot)$ and $\varphi(\cdot)$ denote the Möbius function and Euler's totient function, respectively. For basic properties of these number-theoretic functions, the reader can refer to [3].

Lemma 3.1.2. Let $e>1$ be a positive integer and $q \equiv 1(\bmod e)$ be a prime power. Let $\chi$ be a multiplicative character of order e of $\mathbb{F}_{q}$. For any divisor $w$ of e, let

$$
\begin{equation*}
A_{w}(x)=1+\sum_{k=1}^{e-1} \chi\left(x^{w k}\right) \quad \text { and } \quad A(x)=\sum_{w \mid e} \mu\left(\frac{e}{w}\right) A_{w}(x) \tag{3.1}
\end{equation*}
$$

for $x \in \mathbb{F}_{q}$. Then $A(x)=e$ if $x \in \bigcup_{i \in \mathbb{Z}_{e}^{\times}} C_{i}^{(e)}$ and $A(x)=0$ otherwise.
Proof. First we can reformulate $A_{w}(x)$ in an explicit form, namely

$$
A_{w}(x)= \begin{cases}1, & \text { if } x=0  \tag{3.2}\\ e, & \text { if } x \in \bigcup_{\substack{0 \leq i \leq e-1 \\ e \mid i w}} C_{i}^{(e)}, \\ 0, & \text { otherwise. }\end{cases}
$$

Then the proof can be simply done by the well-known property

$$
\sum_{d \mid n} \mu(d)=\sum_{d \mid n} \mu\left(\frac{n}{d}\right)= \begin{cases}1, & \text { if } n=1 \\ 0, & \text { if } n>1\end{cases}
$$

For $x=0$, we have $A(0)=\sum_{w \mid e} \mu\left(\frac{e}{w}\right)=0$. For a given $x \in C_{i}^{(e)}$, we have $A(x)=\sum_{w|e, e| i w} \mu\left(\frac{e}{w}\right) e=e \sum_{w\left|e, \frac{e}{d}\right| w} \mu\left(\frac{e}{w}\right)$ where $d=\operatorname{gcd}(i, e)$. Suppose $w=$
$\frac{e}{d} \cdot u$, then $\frac{e}{w}=\frac{d}{u}$. We can substitute $e$ and $w$ with $d$ and $u$ in the above sum to get

$$
A(x)=e \sum_{u \mid d} \mu\left(\frac{d}{u}\right)= \begin{cases}e, & \text { if } d=1, \\ 0, & \text { if } d>1,\end{cases}
$$

which is equivalent to saying $A(x)=\left\{\begin{array}{ll}e, & \text { if } x \in C_{i}^{(e)} \\ 0, & \text { otherwise. }\end{array}\right.$ and $i \in \mathbb{Z}_{e}^{\times}$,
Theorem 3.1.3. Let $e \geq 2$ and $t \geq 1$ be positive integers and $q \equiv 1(\bmod e)$ be a prime power. Suppose $a_{j}, b_{j} \in \mathbb{F}_{q}^{*}$ and $c_{j} \in \mathbb{Z}_{e}$ for $1 \leq j \leq t-1$ such that $\left\{a_{j}^{-1} b_{j} \mid 1 \leq j \leq t-1\right\} \cup\{0\}$ is a $t$-subset of $\mathbb{F}_{q}$. Let

$$
\begin{equation*}
X=\left\{x \in \mathbb{F}_{q} \mid x \text { satisfies the following (i) and (ii) }\right\} . \tag{3.3}
\end{equation*}
$$

(i) $x \in C_{i}^{(e)}$ for $i \in \mathbb{Z}_{e}^{\times}$;
(ii) $a_{j} x+b_{j} \in C_{c_{j} i}^{(e)}$ for $1 \leq j \leq t-1$ and $i \in \mathbb{Z}_{e}^{\times}$.

Then $|X|>n$ whenever

$$
\begin{gather*}
q>L(e, t, n):=\left(\frac{c_{1}+\sqrt{c_{1}^{2}+4 \varphi(e) c_{0}}}{2 \varphi(e)}\right)^{2} \text { with }  \tag{3.4}\\
c_{0}:=(e n+t-1) e^{t-1}+e-1 \text { and } c_{1}:=\left(e-w^{*}+\sum_{w \mid e, \mu\left(\frac{e}{w}\right) \neq 0}(e-w)\right) \Psi,
\end{gather*}
$$

where $w^{*}$ is the largest divisor of e with $\mu\left(\frac{e}{w}\right)=-1$ and

$$
\Psi:=\sum_{\ell=1}^{t-1}\binom{t-1}{\ell}(e-1)^{\ell} \ell .
$$

In particular, $X$ is not empty if $q>L(e, t):=L(e, t, 0)$. Furthermore, if $e$ is a prime power of the form $p^{s}$ with $s \geq 1$ and $p$ prime, then $L(e, t, n)=$ $\left(\Psi+\sqrt{\Psi^{2}+c_{0} / \varphi(e)}\right)^{2}$.
Proof. The conditions (i) and (ii) are equivalent to
(i') $x \in \bigcup_{i \in \mathbb{Z}_{e}^{\times}} C_{i}^{(e)}$;
(ii') $x^{e-c_{j}}\left(a_{j} x+b_{j}\right) \in C_{0}^{(e)}$ for $1 \leq j \leq t-1$.
In order to find $L(e, t, n)$, we will define a sum $S$ via a series of character sums of $\mathbb{F}_{q}$. Employing the double counting (estimating) technique on $S$, an inequality with respect to $q$ which guarantees $|X|>n$ will be derived. This is a classical way for showing the asymptotic existence of DFs and "DF-like" structures over $\mathbb{F}_{q}$.

First, let

$$
\begin{equation*}
B(x):=e+\sum_{w \mid e, w \neq e} \mu\left(\frac{e}{w}\right) A_{w}(x), \tag{3.5}
\end{equation*}
$$

where $A_{w}(x)=1+\sum_{k=1}^{e-1} \chi\left(x^{w k}\right)$. It follows from 3.2 that $A_{e}(x)$ is equal to 1 if $x=0$ and is equal to $e$ otherwise. Since $B(x)=e-A_{e}(x)+A(x)$ (where $A(x)$ follows the definition in (3.1) , by Lemma 3.1.2 we have

$$
B(x)= \begin{cases}e-1, & \text { if } x=0  \tag{3.6}\\ e, & \text { if } x \in \bigcup_{i \in \mathbb{Z}_{e}^{\times}} C_{i}^{(e)} \\ 0, & \text { otherwise }\end{cases}
$$

For $1 \leq j \leq t-1$, let $f_{j}(x):=x^{e-c_{j}}\left(a_{j} x+b_{j}\right)$ and

$$
B_{j}(x):=\sum_{k=0}^{e-1} \chi\left(f_{j}^{k}(x)\right)= \begin{cases}1, & \text { if } f_{j}(x)=0  \tag{3.7}\\ e, & \text { if } f_{j}(x) \in C_{0}^{(e)} \\ 0, & \text { otherwise }\end{cases}
$$

Now we consider

$$
\begin{equation*}
S:=\sum_{x \in \mathbb{F}_{q}} B(x) \prod_{j=1}^{t-1} B_{j}(x) \tag{3.8}
\end{equation*}
$$

It is easy to observe that $S=e^{t} m+d$, where $m$ is the number of $x \in \mathbb{F}_{q}$ satisfying conditions ( $\mathrm{i}^{\prime}$ ) and ( $\mathrm{ii}^{\prime}$ ), and $d$ is the contribution when at least one of $x, f_{1}(x), \ldots, f_{t-1}(x)$ is 0 . If $x=0$ and $c_{j} \neq 0$, then $f_{j}(x)=0$ holds for each $1 \leq j \leq t-1$, thus the contribution is $e-1$. Otherwise, if $x \neq 0$ and $f_{j}(x)=0$, then the contribution is at most $e^{t-1}$ for each $1 \leq j \leq t-1$. In order to prove $m \geq n$, it suffices to show $S>e^{t} n+(t-1) e^{t-1}+e-1=c_{0}$.

Now we begin to consider each part of the sum $S=e \sum_{x \in \mathbb{F}_{q}} \prod_{j=1}^{t-1} B_{j}(x)+$ $\sum_{w \mid e, w \neq e} \mu\left(\frac{e}{w}\right) \sum_{x \in \mathbb{F}_{q}} A_{w}(x) \prod_{j=1}^{t-1} B_{j}(x)$. Let

$$
\begin{equation*}
S_{w}:=\sum_{x \in \mathbb{F}_{q}} A_{w}(x) \prod_{j=1}^{t-1} B_{j}(x) \quad \text { and } \quad S_{w}^{-}:=\sum_{x \in \mathbb{F}_{q}}\left(e-A_{w}(x)\right) \prod_{j=1}^{t-1} B_{j}(x) \tag{3.9}
\end{equation*}
$$

Here we simply denote by $\sum^{\prime}$ the sum extended over $\sum_{\ell=1}^{t-1} \sum_{\substack{1 \leq j_{1}<\cdots<j_{\ell} \leq t-1 \\ 1 \leq k_{1}, \ldots, k_{\ell} \leq e-1}}$. Then,

$$
\begin{aligned}
\frac{1}{w} S_{w} & =\sum_{x \in \mathbb{F}_{q}} \sum_{u=0}^{\frac{e}{w}-1} \chi\left(x^{w u}\right) \prod_{j=1}^{t-1} \sum_{k=0}^{e-1} \chi\left(f_{j}^{k}(x)\right) \\
& =\sum_{x \in \mathbb{F}_{q}} 1+U+\sum_{u=1}^{\frac{e}{w}-1} \sum^{\prime} \sum_{x \in \mathbb{F}_{q}} \chi\left(x^{w u} f_{j_{1}}^{k_{1}}(x) \cdots f_{j_{\ell}}^{k_{\ell}}(x)\right)
\end{aligned}
$$

$$
\begin{aligned}
S_{w}^{-}= & \sum_{x \in \mathbb{F}_{q}}\left(e-w-w \sum_{u=1}^{\frac{e}{w}-1} \chi\left(x^{w u}\right)\right) \prod_{j=1}^{t-1} \sum_{k=0}^{e-1} \chi\left(f_{j}^{k}(x)\right) \\
= & \sum_{x \in \mathbb{F}_{q}}(e-w)+(e-w) \sum^{\prime} \sum_{x \in \mathbb{F}_{q}} \chi\left(f_{j_{1}}^{k_{1}}(x) \cdots f_{j_{\ell}}^{k_{\ell}}(x)\right) \\
& -w U-w \sum_{u=1}^{\frac{e}{w}-1} \sum^{\prime} \sum_{x \in \mathbb{F}_{q}} \chi\left(x^{w u} f_{j_{1}}^{k_{1}}(x) \cdots f_{j_{\ell}}^{k_{\ell}}(x)\right)
\end{aligned}
$$

where $\sum_{x \in \mathbb{F}_{q}} 1=q$ and $U=\sum_{u=1}^{\frac{e}{w}-1} \sum_{x \in \mathbb{F}_{q}} \chi\left(x^{w u}\right)=0$. Let

$$
\begin{equation*}
\Gamma_{w}:=S_{w}-w q \quad \text { and } \quad \Gamma_{w}^{-}:=S_{w}^{-}-(e-w) q . \tag{3.10}
\end{equation*}
$$

It follows from Weil's Theorem 1.4.7 that $\left|\sum_{x \in \mathbb{F}_{q}} \chi\left(f_{j_{1}}^{k_{1}}(x) \cdots f_{j_{\ell}}^{k_{\ell}}(x)\right)\right| \leq \ell \sqrt{q}$ and $\left|\sum_{x \in \mathbb{F}_{q}} \chi\left(x^{w u} f_{j_{1}}^{k_{1}}(x) \cdots f_{j_{\ell}}^{k_{\ell}}(x)\right)\right| \leq \ell \sqrt{q}$. We have

$$
\begin{equation*}
\left|\Gamma_{w}\right| \leq w \sum_{u=1}^{\frac{e}{w}-1} \sum^{\prime} \ell \sqrt{q}=(e-w) \Psi \sqrt{q} \quad \text { and } \quad\left|\Gamma_{w}^{-}\right| \leq 2(e-w) \Psi \sqrt{q} \tag{3.11}
\end{equation*}
$$

where $\Psi=\sum^{\prime} \ell=\sum_{\ell=1}^{t-1}\binom{t-1}{\ell}(e-1)^{\ell} \ell$. Now we are in a position to estimate $|S|$. Noting that $\mu\left(\frac{e}{w^{*}}\right)=-1$, we have

$$
S=S_{w^{*}}^{-}+\sum_{\substack{w \mid e \\ w \notin\left\{e, w^{*}\right\}}} \mu\left(\frac{e}{w}\right) S_{w}=\sum_{w \mid e} \mu\left(\frac{e}{w}\right) w q+\Gamma_{w^{*}}^{-}+\sum_{\substack{w \mid e \\ w \notin\left\{e, w^{*}\right\}}} \mu\left(\frac{e}{w}\right) \Gamma_{w},
$$

where $\sum_{w \mid e} \mu\left(\frac{e}{w}\right) w q=\varphi(e) q$. Moreover, we have

$$
\begin{aligned}
\varphi(e) q-|S| \leq|S-\varphi(e) q| & =\left|\Gamma_{w^{*}}^{-}+\sum_{w \mid e, w \notin\left\{e, w^{*}\right\}} \mu\left(\frac{e}{w}\right) \Gamma_{w}\right| \\
& \leq\left|\Gamma_{w^{*}}^{-}\right|+\sum_{w \mid e, w \notin\left\{e, w^{*}\right\}}\left|\mu\left(\frac{e}{w}\right) \Gamma_{w}\right| \\
& \leq\left(e-w^{*}\right) \Psi \sqrt{q}+\sum_{w \mid e, \mu\left(\frac{e}{w}\right) \neq 0}(e-w) \Psi \sqrt{q} \\
& =c_{1} \sqrt{q}
\end{aligned}
$$

where the last inequality follows from (3.11). Obviously, if $q>\left(\frac{c_{1}+\sqrt{c_{1}^{2}+4 \varphi(e) c_{0}}}{2 \varphi(e)}\right)^{2}$, then $|S|>c_{0}$, so that $|X|>n$.

When $e=p^{s}$ is a prime power with $p$ prime, we have $w^{*}=p^{s-1}$ and $c_{1}=2\left(e-p^{s-1}\right) \Psi=2 \varphi(e) \Psi$. Therefore, $L(e, t, n)=\left(\Psi+\sqrt{\Psi^{2}+c_{0} / \varphi(e)}\right)^{2}$.

We want to show that $L(e, t, n)$ is better (smaller) than $Q(e, t, n)$ in most cases. To prove this proposition, we need a simple inequality on integers.

Lemma 3.1.4. Let $n$ be a positive integer and let $d(n)$ denote the number of divisors of $n$. Then $d(n)<\varphi(n)$ whenever $n \notin\{1,2,3,4,6,10,12,30\}$.

Proof. Both $d(\cdot)$ and $\varphi(\cdot)$ are multiplicative functions. Now suppose $p$ is prime and $\alpha \geq 1$. Clearly, $d\left(p^{\alpha}\right)=\alpha+1<\varphi\left(p^{\alpha}\right)=(p-1) p^{\alpha-1}$ whenever $p \geq 5$ and $\alpha \geq 1$, or $p=3$ and $\alpha \geq 2$, or $p=2$ and $\alpha \geq 3$. Furthermore, suppose $p \geq 5$. For $n=2^{2} p^{\alpha}$, or $n=3^{1} p^{\alpha}$, or $n=2^{2} 3^{1} p^{\alpha}, d(n)<\varphi(n)$ always holds. For $n=2 p^{\alpha}$ or $n=2^{1} 3^{1} p^{\alpha}, d(n)<\varphi(n)$ holds except when $p^{\alpha}=5^{1}$. Thus $d(n)<\varphi(n)$ for any positive integer $n \notin\{1,2,3,4,6,10,12,30\}$.

Proposition 3.1.5. For any integers $e \geq 3, e \neq 6, t \geq 2$, and $n \geq 0, L(e, t, n)<$ $Q(e, t, n)$. In particular, if $t \geq 3$ and $e=p^{s}$ with $p$ prime and $s \geq 1$, then $\frac{(e-1)^{2}}{4} L(e, t)<Q(e, t)$.

Proof. For $e \notin\{3,4,6,10,12,30\}$, it follows from Lemma3.1.4 that $d(e)<\varphi(e)$, thus $\gamma(e):=e-w^{*}+\sum_{w \mid e, \mu\left(\frac{e}{w}\right) \neq 0}(e-w) \leq e-w^{*}+\sum_{w \mid e, w \neq e}(e-w) \leq$ $\sum_{w \mid e}(e-1)=(e-1) d(e)<(e-1) \varphi(e)$. For $e \in\{3,4,10,12,30\}$, it can be directly verified that $\gamma(e) \leq(e-1) \varphi(e)$. Moreover,

$$
\begin{aligned}
\frac{U}{e-1}-\Psi & =\sum_{\ell=1}^{t-1}\binom{t}{\ell+1}(e-1)^{\ell} \ell-\sum_{\ell=1}^{t-1}\binom{t-1}{\ell}(e-1)^{\ell} \ell \\
& =\sum_{\ell=1}^{t-1} \frac{t-\ell-1}{\ell+1}\binom{t-1}{\ell}(e-1)^{\ell} \ell
\end{aligned}
$$

which is greater than or equal to $e-1$ when $t \geq 3$, and is equal to 0 if $t=$ 2. Hence, $\frac{c_{1}}{\varphi(e)}=\frac{\gamma(e)}{\varphi(e)} \Psi \leq \frac{(e-1) \varphi(e)}{\varphi(e)} \cdot \frac{U}{e-1}=U$. Furthermore, since $\frac{c_{0}}{\varphi(e)}<$ $c_{0}=(e n+t-1) e^{t-1}+e-1<(e n+t) e^{t-1}$, we have $\frac{c_{1}+\sqrt{c_{1}^{2}+4 \varphi(e) c_{0}}}{\varphi(e)}<U+$ $\sqrt{U^{2}+4 t e^{t-1}}$, which immediately implies $L(e, t)<Q(e, t)$.

Next, in the case of $n=0$, we suppose $t \geq 3$ and $e$ is a prime power. Proving $\frac{(e-1)^{2}}{4} L(e, t)<Q(e, t)$ is equivalent to showing $\Psi+\sqrt{\Psi^{2}+c_{0} / \varphi(e)}<$ $\frac{1}{e-1}\left(U+\sqrt{U^{2}+4 t e^{t-1}}\right)$. As shown above, $\Psi<\frac{U}{e-1}$. So it suffices to show $\Psi^{2}+\frac{c_{0}}{\varphi(e)}<\frac{U^{2}}{(e-1)^{2}}+\frac{4 t e^{t-1}}{(e-1)^{2}}$. This inequality can be obtained by combining $\left(\frac{U}{e-1}\right)^{2}-\Psi^{2}>\frac{e-1}{2} t e^{t-1} \geq t e^{t-1}$ and $\frac{c_{0}}{\varphi(e)}-\frac{4 t e^{t-1}}{(e-1)^{2}}<\frac{c_{0}}{\varphi(e)}=\frac{(t-1) e^{t-1}+e-1}{\varphi(e)}<$ $\frac{t e^{t-1}}{\varphi(e)}<t e^{t-1}$, where the first inequality follows from $\frac{U}{e-1}-\Psi \geq e-1$ and $\frac{U}{e-1}+\Psi=\sum_{\ell=1}^{t-1}\binom{t-1}{\ell}(e-1)^{\ell}\left(\frac{\ell t}{\ell+1}+\ell\right)>\sum_{\ell=1}^{t-1}\binom{t-1}{\ell}(e-1)^{\ell}\left(\frac{t}{2}+1\right)=\left(\frac{t}{2}+\right.$ 1) $\left(e^{t-1}-1\right)>\frac{t}{2} e^{t-1}$.

### 3.2 Direct constructions and asymptotic existence of grid-block difference families

First, we need some notation. Let $G$ be an additive group and let $A=$ $\left\{a_{1}, a_{2}, \ldots, a_{t}\right\}$ be a subset of $G$. Let $\Delta A=\left\{a_{i}-a_{j} \mid 1 \leq i, j \leq t, i \neq j\right\}$ and $\Delta^{+} A=\left\{a_{i}-a_{j} \mid 1 \leq i<j \leq t\right\}$. Further, suppose $\mathrm{B}=\left[b_{i j}\right]_{r \times k}$ is an $r \times k$ grid-block whose elements are in $G$. Similarly, let

$$
\Delta \mathrm{B}=\left(\bigcup_{i=1}^{r} \Delta\left\{b_{i 1}, b_{i 2}, \ldots, b_{i k}\right\}\right) \cup\left(\bigcup_{j=1}^{k} \Delta\left\{b_{1 j}, b_{2 j}, \ldots, b_{r j}\right\}\right)
$$

Let $\mathcal{B}$ be a collection of grid-blocks, and denote $\Delta \mathcal{B}=\bigcup_{\mathrm{B} \in \mathcal{B}} \Delta \mathrm{B}$.
Lemma 3.2.1 (131] Lemma 1.7). Let $e=r k(r+k-2) / 2$ and let $q \equiv 1$ $(\bmod 2 e)$ be a prime power. If there exists an $r \times k$ array A over $\mathbb{F}_{q}$ such that $\Delta^{+} \mathrm{A}$ is a system of representatives of $\mathcal{C}^{(e)}$, then there exists a $\left(q, L_{r \times k}, 1\right)-D F$ over $\mathbb{F}_{q}$.

Next, we propose more theorems which extend Lemma 3.2.1.

### 3.2.1 Grid-block difference families with a multiplier of order 3

In this section, we give a direct construction of a $\left(q, L_{r, 3 u}, 1\right)$-DF with a multiplier of order 3 , which requires $q \equiv 1(\bmod 3 \operatorname{ur}(r+3 u-2))$.

Theorem 3.2.2. For any positive integer $r$ and $u$, let $e=\frac{r u(r+3 u-2)}{2}$ and let $q \equiv 1(\bmod 6 e)$ be a prime power. Let $g$ be a primitive element in $\mathbb{F}_{q}$ and $\omega$ be a primitive cubic root of unity of $\mathbb{F}_{q}$. Suppose there exist $x_{1}, \ldots, x_{r-1}, y_{1}, \ldots, y_{u-1} \in$ $\mathbb{F}_{q}^{*}$ such that $\frac{1}{\omega-1} H_{r, u}$ forms a complete system of representatives of $\mathcal{C}^{(e)}$, where

$$
H_{r, u}=(\omega-1) \cdot X \cdot Y \cup X \cdot \Delta_{\omega}^{+} Y \cup Y \cdot \Delta^{+} X
$$

with

$$
\begin{aligned}
\Delta_{\omega}^{+} Y & =\left\{y_{j_{1}}-y_{j_{2}} \omega^{k} \mid 0 \leq j_{1}<j_{2} \leq u-1,0 \leq k \leq 2\right\} \\
X & =\left\{x_{0}, x_{1}, \ldots, x_{r-1}\right\} \\
Y & =\left\{y_{0}, y_{1}, \ldots, y_{u-1}\right\}
\end{aligned}
$$

and $x_{0}=y_{0}=1$. Let $\mathrm{B}=\left[b_{i j}\right]_{r \times 3 u}$ be a grid-block with $\left.b_{i j}=x_{i-1} y_{\left\lceil\frac{j-1}{3}\right.}\right\rceil^{\omega^{j-1}}$ for $1 \leq i \leq r$ and $1 \leq j \leq 3 u$. Then $\mathcal{B}=\left\{g^{e i} \mathrm{~B} \left\lvert\, 0 \leq i<\frac{q-1}{6 e}\right.\right\}$ forms $a$ $\left(q, L_{r, 3 u}, 1\right)-D F$ in $\mathbb{F}_{q}$.
Proof. It is sufficient to show that $\Delta \mathcal{B}=\mathbb{F}_{q}^{*}$. Let $n=\frac{q-1}{6 e}$. First, we have $\left.\left.\Delta \mathrm{B}=X \cdot \Delta\left(\left\{1, \omega, \omega^{2}\right\} \cdot Y\right\}\right) \cup\left(\left\{1, \omega, \omega^{2}\right\} \cdot Y\right\}\right) \cdot \Delta X=\{1,-1\} \cdot\left\{1, \omega, \omega^{2}\right\} \cdot H_{r, u}=$ $C^{(e n)} \cdot H_{r, u}$ with $\left|H_{r, u}\right|=e$. Note that $T:=\left\{g^{e i} \mid 0 \leq i<n\right\}$ is a system of representatives for the cosets of $C^{(e n)}$ in $C^{(e)}$. With the assumption that $H_{r, u}$ forms a complete system of representatives of $\mathcal{C}^{(e)}$, we have $\Delta \mathcal{B}=T \cdot \Delta \mathrm{~B}=$ $C^{(e)} \cdot H_{r, u}=\mathbb{F}_{q}^{*}$.

Table 3.1: $\left(p, L_{6,2}, 1\right)$-DF that cannot obtained from Theorem 3.2.2

| $p$ | $g$ | $\omega$ | $x$ | $y$ | $z$ | $p$ | $g$ | $\omega$ | $x$ | $y$ | $z$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 73 | 5 | 8 | 43 | 59 | 2 | 433 | 5 | 198 | 116 | 125 | 114 |
| 577 | 5 | 363 | 433 | 5 | 80 | 937 | 5 | 614 | 383 | 5 | 417 |
| 1009 | 11 | 374 | 385 | 322 | 981 | 1153 | 5 | 650 | 101 | 5 | 819 |
| 1297 | 17 | 931 | 95 | 513 | 113 | 1657 | 11 | 70 | 1258 | 11 | 1232 |
| 2089 | 7 | 1262 | 1266 | 49 | 515 | 3313 | 11 | 2189 | 939 | 121 | 1388 |
| 3529 | 17 | 3080 | 795 | 289 | 2530 | 7489 | 7 | 5021 | 1616 | 343 | 7176 |

Example 3.2.3. Let $r=2, u=1$, and $q=37$. Take $g=2$, then $\omega=g^{12}=26$.
Take $x_{1}=2$, we have $\left(\zeta-1, x_{1}(\zeta-1), x_{1}-1\right)=(25,13,1) \in C_{1}^{(3)} \times C_{2}^{(3)} \times C_{0}^{(3)}$. Then,

$$
\left\{\left[\begin{array}{ccc}
1 & 26 & 10 \\
2 & 15 & 20
\end{array}\right],\left[\begin{array}{ccc}
6 & 8 & 23 \\
12 & 16 & 9
\end{array}\right]\right\}
$$

forms a cyclic $\left(37, L_{2,3}, 1\right)$-DF.
Remark. When $u=2$ and $r=1$, Theorem 3.2 .2 becomes Wilson 120 Theorem 11 for a $(q, 6,1)$-DF (see also Chen and Zhu [28]).
Remark. For $r=u=2$, we have checked that Theorem 3.2 .2 can be applied to every prime $p \equiv 1(\bmod 72)$ with $p<10^{7}$ except when $p \in P$, where $P=\{p \equiv 1$ $(\bmod 72)$ is prime $\mid p \leq 1657\} \cup\{2089,3313,3529,7489\}$. In other words, $H_{2,2}$ never forms a system of representatives of $\mathcal{C}^{(12)}$ for any choice of a pair $\left(x_{1}, y_{1}\right)$ in $\mathbb{F}_{p}^{*}$ with $p \in P$. However, by introducing an extra variable, a ( $p, L_{2,6}, 1$ )-DF with a multiplier of order 3 can be constructed for any $p \in P$. With the parameters listed in Table 3.1, $\mathcal{B}:=\left\{g^{12 i} \mathrm{~B} \left\lvert\, 0 \leq i<\frac{p-1}{72}\right.\right\}$ forms a $\left(p, L_{2,6}, 1\right)$-DF, where

$$
\mathrm{B}=\left[\begin{array}{cccccc}
1 & \omega & \omega^{2} & y & \omega y & \omega^{2} y \\
x & \omega x & \omega^{2} x & z & \omega z & \omega^{2} z
\end{array}\right]
$$

In addition, other examples of $\left(p, L_{2,6}, 1\right)$-DFs with $p \in\{73,433\}$ can be found from Wang and Colbourn 115 .

Theorem 3.2.4. For any positive integer $r$ and $u$, let $e=\frac{r u(r+3 u-2)}{2}$. Then $a\left(q, L_{r, 3 u}, 1\right)-D F$ exists for all $q>L(e, \max \{r, 3 u-2\})$ with $q \equiv{ }_{1}^{2}(\bmod 6 e) a$ prime power.

Proof. With the notation of $X, Y$, and $\Delta_{\omega}^{+} Y$ in Theorem 3.2.2, it suffices to find $X, Y \subset \mathbb{F}_{q}^{*}$ satisfying the following conditions in two cases:
(a) If $r$ is odd,
(a1) $x_{i} \in C_{i}^{(e)}$ for $i \in[r-1]$;
(a2) $y_{j} \in C_{j r}^{(e)}$ for $j \in[u-1]$;
(a3) $\frac{1}{\omega-1} \Delta_{\omega}^{+} Y$ forms a system of representatives of $\left\{C_{u r+j r}^{(e)} \left\lvert\, j \in\left[0, \frac{3 u(u-1)}{2}-\right.\right.\right.$
(a4) $\frac{1}{\omega-1} \Delta^{+} X$ forms a system of representatives of $\bigcup_{\ell=0}^{\frac{r-3}{2}}\left\{C_{s_{\ell}+i}^{(e)} \mid i \in[0, r-\right.$ $1]\}$ with $s_{\ell}=\frac{r u(3 u-1+2 \ell)}{2}$.
(b) If $r$ is even,
(b1) $x_{i} \in C_{i}^{(e)}$ for $i \in\left[\frac{r}{2}-1\right]$ and $x_{i} \in C_{\frac{r}{2} \cdot \frac{u(3 u-1)}{2}+i}^{(e)}$ for $i \in\left[\frac{r}{2}, r-1\right]$;
(b2) $y_{j} \in C_{j \frac{r}{2}}^{(e)}$ for $j \in[u-1]$;
(b3) $\frac{1}{\omega-1} \Delta_{\omega}^{+} Y$ forms a system of representatives of $\left\{\left.C_{u \frac{r}{2}+j \frac{r}{2}}^{(e)} \right\rvert\, j \in\left[0, \frac{3 u(u-1)}{2}-\right.\right.$ 1] $\} ;$
(b4) $\frac{1}{\omega-1} \Delta^{+} X$ forms a system of representatives of $\bigcup_{\ell=0}^{r-2}\left\{C_{s_{\ell}+i}^{(e)} \left\lvert\, i \in\left[0, \frac{r}{2}-\right.\right.\right.$ $1]\}$ with $s_{\ell}=\frac{r u(3 u-1+\ell)}{2}$.

The subscript $i$ of $C_{i}^{(e)}$ can be regarded as an element in $\mathbb{Z}_{e}$. Clearly, by multiplying all the above subscripts by any unit in $\mathbb{Z}_{e}^{\times}$, we can obtain another quadruple of conditions which is also admissible. Note that there are precisely $r$ (resp. $3 u-2$ ) conditions with respect to $x_{i}$ (resp. $y_{j}$ ) for each $i \in[r-1]$ (resp. $j \in[u-1]$ ). By repeatedly employing Theorem 3.2 .2 for each $x_{i}$ and $y_{j}$, it can be guaranteed that $X, Y \subset \mathbb{F}_{q}^{*}$ satisfying the above conditions exist whenever $q>L\left(\frac{e}{6}, \max \{r, 3 u-2\}\right)$.

In Table 3.2, some existence bounds for a $\left(q, L_{r, 3 u}, 1\right) \mathrm{DF}$ are listed, where $L(e, t)=L(e, t, 0)$ and $Q(e, t)=Q(e, t, 0)$ are defined in Theorems 3.1.3 and 1.4.8 respectively. It is remarkable that a $\left(q, L_{1,6}, 1\right) \mathrm{DF}$ is nothing more than a $(q, 6,1)-\mathrm{DF}$, and the value $L(5,4)$ was first derived by Chen and Zhu [28] via the same estimation. The last columns for $\left(q, L_{1,9}, 1\right)$-DF and $\left(q, L_{1,12}, 1\right)$-DF show the known bound obtained by Chen and Zhu 30 without considering any multiplier. In addition, it can be observed that $L(e, t)$ is more effective when $e$ is a prime power. For real world applications to biology experiments, $8 \times 12$ gridblock designs are the most important. Here we also give a bound for $\left(q, L_{8,12}, 1\right)$ DF , although it is still quite large and unsatisfactory.

### 3.2.2 Row-radical grid-block difference families

As an analogue of radical DFs, it is natural to consider row-radical (columnradical) $\left(q, L_{r, k}, 1\right)$-DFs.

Definition 3.2.5 (row-radical DF). An elementary abelian $\left(q, L_{r, k}, 1\right)$-DF is row-radical (resp., column-radical) if the rows (resp., columns) are cosets of $C^{\left(\frac{q-1}{k}\right)}$ (resp., $\left.C^{\left(\frac{q-1}{r}\right)}\right)$ in all base grid-blocks. Moreover, an elementary abelian ( $q, L_{r, k}, 1$ )-DF is radical if it is both row-radical and column-radical.

Table 3.2: Improved existence bounds for $\left(q, L_{r, 3 u}, 1\right)$-DFs

| $r \times 3 u$ | $t$ | $e$ | $L(e, t)$ | $Q(e, t)$ | Remark / Reference |
| :--- | :--- | :--- | :--- | :--- | :--- |
| $4 \times 3$ | 4 | 10 | $3.3215 \times 10^{8}$ | $6.7606 \times 10^{8}$ |  |
| $5 \times 3$ | 5 | 15 | $1.1810 \times 10^{12}$ | $7.7528 \times 10^{12}$ |  |
| $1 \times 6$ | 4 | 5 | $3.6019 \times 10^{5}$ | $1.8944 \times 10^{6}$ | $=L(5,4)[28]$ |
| $2 \times 6$ | 4 | 12 | $1.2702 \times 10^{9}$ | $3.0578 \times 10^{9}$ |  |
| $3 \times 6$ | 4 | 21 | $2.1179 \times 10^{10}$ | $2.9855 \times 10^{11}$ |  |
| $4 \times 6$ | 4 | 32 | $3.6277 \times 10^{10}$ | $9.0882 \times 10^{12}$ |  |
| $5 \times 6$ | 5 | 45 | $8.5043 \times 10^{15}$ | $5.1496 \times 10^{17}$ |  |
| $1 \times 9$ | 7 | 12 | $1.5171 \times 10^{16}$ | $3.7671 \times 10^{16}$ | $4.7864 \times 10^{20} \quad[30$ |
| $2 \times 9$ | 7 | 27 | $2.0042 \times 10^{19}$ | $3.6060 \times 10^{21}$ |  |
| $1 \times 12$ | 10 | 22 | $4.2694 \times 10^{27}$ | $5.1514 \times 10^{28}$ | $4.1773 \times 10^{31}[30$ |
| $8 \times 12$ | 10 | 288 | $8.4072 \times 10^{47}$ | $1.2387 \times 10^{51}$ |  |

Note that, for the radical cases, we necessarily have $\operatorname{gcd}(r, k)=\operatorname{gcd}(r k, 2)=$ 1 , and the vertex-set of each base grid-block is a coset of $C^{\left(\frac{q-1}{r k}\right)}$. It is clear that a row-radical (column-radical) grid-block DF must be a disjoint grid-block DF.

Now, we concentrate on direct constructions of elementary abelian ( $q, L_{r, k}, 1$ )DDF with odd $k$. First, we consider a row-radical $\left(q, L_{r, k}, 1\right)$-DF over $\mathbb{F}_{q}$, which necessarily requires $k$ to be odd and $q \equiv 1(\bmod r k(r+k-2))$. The following is the main theorem of our construction, which can imply a series of existence results.

Theorem 3.2.6 (Row-radical DF). Let $k$ be an odd integer. Suppose $q=$ $n \cdot r k(r+k-2)+1$ is a prime power for a positive integer $n$. There exists a row-radical $\left(q, L_{r, k}, 1\right)$-DF over $\mathbb{F}_{q}$ if there exist $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{r-1} \in \mathbb{F}_{q}^{*}$ such that $\frac{1}{\zeta_{k}-1} X_{r, k}$ forms a system of representatives of the cosets of $C^{\left(f^{e+1}\right)}$ in $C^{\left(f^{e}\right)}$ for a suitable e such that $f^{e}$ divides $n$, where

$$
\begin{aligned}
A_{r} & =\left\{1, \alpha_{1}, \alpha_{2}, \ldots, \alpha_{r-1}\right\} \\
H_{k} & =\left\{\zeta_{k}-1, \zeta_{k}^{2}-1, \ldots, \zeta_{k}^{\frac{k-1}{2}}-1\right\} \\
X_{r, k} & =A_{r} \cdot H_{k} \cup \Delta^{+} A_{r}
\end{aligned}
$$

and $f=\left|X_{r, k}\right|=\frac{r}{2}(r+k-2)$.
Proof. Let $Z=\left\{1, \zeta_{k}, \ldots, \zeta_{k}^{k-1}\right\}=C^{(2 n f)}$. Since $-\zeta_{k}^{\frac{k-1}{2}+j}\left(\zeta_{k}^{\frac{k-1}{2}+1-j}-1\right)=$ $\zeta_{k}^{\frac{k-1}{2}+j}-1$, we have $-\left(\zeta_{k}^{\frac{k-1}{2}+1-j}-1\right) Z=\left(\zeta_{k}^{\frac{k-1}{2}+j}-1\right) Z$ for any $j \in\left[\frac{k-1}{2}\right]$. Hence, $\Delta Z=\left\{\zeta_{k}^{i}\left(\zeta_{k}^{j}-1\right) \mid i, j \in \mathbb{Z}_{k}, j \neq 0\right\}=Z \cdot\left\{\zeta_{k}^{j}-1, \left.\zeta_{k}^{\frac{k-1}{2}+j}-1 \right\rvert\, 1 \leq j \leq\right.$ $\left.\frac{k-1}{2}\right\}= \pm Z \cdot H_{k}=C^{(n f)} \cdot H_{k}$.

Let $\mathrm{B}=\left[b_{i j}\right]_{r \times k}$ with $b_{i j}=\alpha_{i-1} \zeta_{k}^{j-1}$ for $i \in[r]$ and $j \in[k]$, where $\alpha_{0}=1$.

Then $\Delta \mathrm{B}=A_{r} \cdot \Delta Z \cup Z \cdot \Delta A_{r}=C^{(n f)} \cdot X_{r, k}$. Let

$$
\mathcal{B}=\left\{g^{f^{e+1} i+j} \cdot \mathrm{~B} \left\lvert\, i \in\left[\frac{n}{f^{e}}\right]\right., j \in\left[f^{e}\right]\right\} .
$$

Then, $\Delta \mathcal{B}=\left\{g^{f^{e+1} i+j} \cdot X_{r, k} \left\lvert\, i \in\left[\frac{n}{f^{e}}\right]\right., j \in\left[f^{e}\right]\right\} \cdot C^{(n f)}$. Now we show the distinctness of all the elements in $\frac{1}{\zeta_{k}-1} \Delta \mathcal{B}$. Assume there exist $i_{1}, i_{2} \in\left[\frac{n}{f^{e}}\right], j_{1}, j_{2} \in$ $\left[f^{e}\right], x_{1}, x_{2} \in \frac{1}{\zeta_{k}-1} X_{r, k}$, and $u_{1}, u_{2} \in\left[\frac{q-1}{n f}\right]$ with $\left(i_{1}, j_{1}, x_{1}, u_{1}\right) \neq\left(i_{2}, j_{2}, x_{2}, u_{2}\right)$ such that $g^{f^{e+1} i_{1}+j_{1}} \cdot x_{1} \cdot g^{n f u_{1}}=g^{f^{e+1} i_{2}+j_{2}} \cdot x_{2} \cdot g^{n f u_{2}}$, where $x_{1}=g^{f^{e+1} s_{1}+f^{e} t_{1}}$ and $x_{2}=g^{f^{e+1} s_{2}+f^{e} t_{2}}$. This is equivalent to saying $f^{e+1} i_{1}+j_{1}+f^{e+1} s_{1}+f^{e} t_{1}+$ $n f u_{1} \equiv f^{e+1} i_{2}+j_{2}+f^{e+1} s_{2}+f^{e} t_{2}+n f u_{2}(\bmod q-1)$ which implies $j_{1} \equiv j_{2}$ $\left(\bmod f^{e}\right)$ and $\left(i_{1}+s_{1}+\frac{n}{f^{e}} u_{1}\right) f+t_{1} \equiv\left(i_{2}+s_{2}+\frac{n}{f^{e}} u_{2}\right) f+t_{2}\left(\bmod \frac{q-1}{f^{e}}\right)$. The second congruence implies that $t_{1} \equiv t_{2}\left(\bmod \frac{q-1}{f^{e}}\right)$. Recalling that $\frac{1}{\zeta_{k}-1} X_{r, k}$ forms a system of representatives of the cosets of $C^{\left(f^{e+1}\right)}$ in $C^{\left(f^{e}\right)}$, there must be $x_{1}=x_{2}$ as well. So, $s_{1}=s_{2}$. Hence, $i_{1}+\frac{n}{f^{e}} u_{1} \equiv i_{2}+\frac{n}{f^{e}} u_{2}\left(\bmod \frac{q-1}{f^{e+1}}\right)$, i.e., $i_{1} \equiv i_{2}\left(\bmod \frac{q-1}{f^{e+1}}\right)$ and $u_{1} \equiv u_{2}\left(\bmod \frac{q-1}{n f}\right)$. This conclusion contradicts the assumption $\left(i_{1}, j_{1}, x_{1}, u_{1}\right) \neq\left(i_{2}, j_{2}, x_{2}, u_{2}\right)$. In summary, $\left|\frac{1}{\zeta_{k}-1} \Delta \mathcal{B}\right|=q-1$. Therefore, $\Delta \mathcal{B}=\left(\zeta_{k}-1\right) \cdot \mathbb{F}_{q}^{*}=\mathbb{F}_{q}^{*}$.

For the case when $r=1$, we can set $A_{r}=\{1\}$ and $\Delta^{+} A_{r}=\emptyset$. Then Theorem 3.2.6 becomes the construction of radical DF due to Bose [10] ( $e=0$ and $k \in\{3,5\}$ ), Wilson [120] ( $e=0$ and $k$ odd), and Buratti [15, 16] (who also showed the necessity when $k \leq 7$ ).

For $r=2$ and $k=3$, by using Theorem 1.4.8 we can meet the criteria in Theorem 3.2 .6 with $e=0$ for any admissible $q>26$. Here we only state the conclusion without proof, since the existence has been presented in 9.

Corollary 3.2.7. There exists an elementary abelian $\left(q, L_{2,3}, 1\right)$-DDF for any prime power $q \equiv 1(\bmod 18)$.

Actually, when $r=2$, we have $X_{2, k}=\left\{1, \alpha_{1}\right\} \cdot H_{k} \cup\left\{\alpha_{1}-1\right\}$ with cardinality $\left|X_{2, k}\right|=k$. For the case when $p$ is a prime and $e=0$ in Theorem3.2.6 in order for $X_{2, k}$ to form a system of representatives of $\mathcal{C}^{(k)}, p$ is necessarily a good prime (cf. Definition 3.3.1). In particular, for $r=2$ and $k=5$, the cyclotomic condition $1+\zeta_{5} \notin C^{(5)}$ for $p$ to be good is sufficient when $p$ is an admissible prime.

First we present the equivalence between the existence of a row-radical $\left(q, L_{r, k}, 1\right)$-DF and a row-radical $\left(q^{m}, L_{r, k}, 1\right)$-DF.

Lemma 3.2.8. Let $q \equiv 1(\bmod r k(r+k-2))$ be a prime power. Theorem 3.2.6 gives a $\left(q, L_{r, k}, 1\right)$-DDF if and only if it gives a $\left(q^{m}, L_{r, k}, 1\right)$-DDF for any positive integer $m$.

Proof. A $\left(q, L_{r, k}, 1\right)$-DDF obtained from Theorem 3.2 .6 must be of the form $\mathcal{B}=\{x \mathrm{~B} \mid x \in X\}$, where $X \subset \mathbb{F}_{q}^{*}$. Let $T$ be a transversal (a complete system of representatives) of the cosets of $\mathbb{F}_{q}^{*}$ in $\mathbb{F}_{q^{m}}^{*}$. Since $q^{m}-1$ can also be divided
by $r k(r+k-2)$, then $\mathcal{B}_{m}:=\{t x \mathrm{~B} \mid x \in X, t \in T\}$ forms a $\left(q^{m}, L_{r, k}, 1\right)$-DDF in $\mathbb{F}_{p^{m}}$. It is clear that each row of any grid-block in $\mathcal{B}_{m}$ forms a coset of the $k$ th root of unity in $\mathbb{F}_{q^{m}}$, thus $\mathcal{B}_{m}$ can be derived by Theorem 3.2.6.

Conversely, suppose $\mathcal{B}_{m}$ consists of the base grid-blocks of a $\left(q^{m}, L_{r, k}, 1\right)$ DDF obtained from Theorem 3.2.6. Then $\mathcal{B}_{m}$ must be of the form $\{y \mathrm{~B} \mid y \in Y\}$ for some $Y \subset \mathbb{F}_{q^{m}}^{*}$. It is easily seen that $\mathcal{B}^{\prime}:=\left\{x \mathrm{~B} \mid x \in Y \cap \mathbb{F}_{q}\right\}$ forms a $\left(q, L_{r, k}, 1\right)$-DDF in $\mathbb{F}_{q}$. Moreover, each row of any grid-block in $\mathcal{B}^{\prime}$ is a coset of the $k$ th root of unity in $\mathbb{F}_{q}$. Hence $\mathcal{B}_{m}$ can be derived by Theorem 3.2.6.

As a generalization of Corollary 3.2.7 on a $\left(q, L_{2,3}, 1\right)$-DDF, we consider a row-radical $\left(q, L_{r, 3}, 1\right)$-DF over $\mathbb{F}_{q}$ with $r \geq 3$. It is necessary to suppose $v \equiv 1$ $(\bmod 3 r(r+1))$.

Theorem 3.2.9 (Asymptotic existence for $k=3$ ). For any positive integer $r \geq 3$, let $q \equiv 1\left(\bmod 3 r(r+1)\right.$ ) be a prime power with $q>r^{2}\binom{r+1}{2}^{2 r}$. Then there exists a row-radical $\left(q, L_{r, 3}, 1\right)$-DF over $\mathbb{F}_{q}$.

Proof. Set $k=3$ and $e=0$ in Theorem 3.2.6. We have $X_{r, 3}=\left(1-\zeta_{3}\right) A_{r} \cup \Delta^{+} A_{r}$ with $\left|X_{r, 3}\right|=r+\binom{r}{2}=\binom{r+1}{2}$, which we want to form a system of representatives of $\mathcal{C}^{\left(\left|X_{r, 3}\right|\right)}$. Then it suffices to bound $q$ by setting $n=\left|X_{r, 3}\right|=\binom{r+1}{2}$ and $s=r$ in Theorem 1.4.8 and using the fact $s^{2} n^{2 s}>Q(n, s)$ to meet the claim.

Corollary 3.2.10. There exists an elementary abelian ( $q, L_{3,3}, 1$ )-DDF for any prime power $q \equiv 1(\bmod 36)$.

Proof. It follows from Theorem 3.2 .9 that a $\left(q, L_{3,3}, 1\right)$-DDF exists for any prime power $q \equiv 1(\bmod 36)$ with $q \geq 105841>Q\left(\binom{4}{2}, 3\right)$. Let $S_{1}=\{p:$ prime $\mid$ $p \equiv 1(\bmod 36), p<105841\}$. For non-prime $q$, by Lemma 3.2 .8 , we only need to check the prime powers in $S_{2}:=\left\{p^{2} \mid p \equiv-1, \pm 17(\bmod 36), p \leq 307, p:\right.$ prime $\} \cup\left\{13^{3}, 5^{6}\right\}$. We have individually checked every $q \in S_{1} \cup S_{2}$ and succeeded in constructing a desired $\left(q, L_{3,3}, 1\right)$-DF by Theorem 3.2.6 with $e=0$.

Similarly, we can consider a row-radical ( $p, L_{r, 5}, 1$ )-DF.
Theorem 3.2.11 (Asymptotic existence for $k=5$ ). For any positive integer $r \geq 3$, let $p \equiv 1(\bmod 5 r(r+3))$ be a prime with $p \geq r^{2 r+2}\left(\frac{r+3}{2}\right)^{2 r}$ such that $\zeta_{5}+1 \notin C^{\left(\frac{r(r+3)}{2}\right)}$. Then there exists a row-radical $\left(p, L_{r, 5}, 1\right)-D F$.

Proof. This is similar to the proof of Theorem 3.2.9. In this case, $X_{r, 5}=$ $\left\{1-\zeta_{5}, 1-\zeta_{5}^{2}\right\} \cdot A_{r} \cup \Delta A_{r}$ is desired to form a system of representatives of $\mathcal{C}{ }^{\left(\left|X_{r, 5}\right|\right)}$. When $\zeta_{5}+1 \notin C^{\left(\left|X_{r, 5}\right|\right)}$, it suffices to bound $p$ by setting $n=\left|X_{r, 5}\right|=\frac{r(r+3)}{2}$ and $s=r$ in Theorem 1.4.8 and use the fact $s^{2} n^{2 s}>Q(n, s)$ to meet the claim.

For $k \geq 5$, we will discuss row-radical $\left(p, L_{2, k}, 1\right)$-DFs in details in Section 3.3. For row-radical $\left(p, L_{3,5}, 1\right)$-DFs, we have the following:

Corollary 3.2.12. There exists a row-radical ( $p, L_{3,5}, 1$ )-DF for any prime $p \equiv$ $1(\bmod 90)$ satisfying $\zeta_{5}+1 \notin C^{(9)}$.

Proof. By setting $r=3$ in Theorem 3.2 .11 , we can obtain a row-radical $\left(p, L_{3,5}, 1\right)$ DDF satisfying the claimed condition with $p>1,479,142>Q(9,3)$. We individually checked the remaining admissible primes satisfying $\zeta_{5}+1 \notin C^{(9)}$ and succeeded in constructing the difference families.

Next, we consider a radical (namely, both row- and column-radical) ( $q, L_{r, k}, 1$ )DF. This is actually a very special but nice case of Theorem 3.2 .6 obtained by setting $A_{r}=\left\{1, \zeta_{r}, \zeta_{r}^{2}, \ldots, \zeta_{r}^{r-1}\right\}$. In this case, we can obtain a sufficient cyclotomic condition for the existence. We indicate such a DF is nice because it is quite simple to calculate, in particular when $q$ is not so huge.

Theorem 3.2.13 (Radical DF). Let $r>k>1$ be odd integers which are relatively prime. Suppose $q=n \cdot r k(r+k-2)+1$ is a prime power for a positive integer $n$. There exists a radical $\left(q, L_{r, k}, 1\right)$-DF over $\mathbb{F}_{q}$ if $\frac{1}{\zeta_{k}-1} Y_{r, k}$ forms a system of representatives of the cosets of $C^{\left(h^{e+1}\right)}$ in $C^{\left(h^{e}\right)}$ for a suitable e such that $h^{e}$ divides $n$, where

$$
\begin{aligned}
Y_{r, k} & =H_{r} \cup H_{k} \\
& =\left\{\zeta_{r}-1, \zeta_{r}^{2}-1, \ldots, \zeta_{r}^{\frac{r-1}{2}}-1\right\} \cup\left\{\zeta_{k}-1, \zeta_{k}^{2}-1, \ldots, \zeta_{k}^{\frac{k-1}{2}}-1\right\}
\end{aligned}
$$

and $h=\left|Y_{r, k}\right|=\frac{r+k-2}{2}$.
Proof. Since $\Delta^{+} A_{r}=A_{r} \cdot H_{r}$, we have $X_{r, k}=A_{r} \cdot\left(H_{k} \cup H_{r}\right)=C^{(n \cdot k(r+k-2))} \cdot Y_{r, k}$. Then it is easy to deduce the conclusion from Theorem 3.2.6 by using the fact that $\operatorname{gcd}(r, k)=1$.

Remark. Among all the 4679 admissible primes with $p<1,479,141$ (the bound obtained from Theorem 3.2.11), there are 1058 of them satisfying the existence condition for a radical DF (including 1020 primes with $e=0,36$ primes with $e=1$, and two primes with $e=2$ which are also the only two satisfying $\zeta_{5}+1 \in C^{(9)}$ ), where the first two primes are $p=541$ (with $e=1$ ) and $p=1171$ (with $e=0$ ). Moreover, besides the "radical" ones, only 500 of the remaining admissible primes lead to $\zeta_{5}+1 \in C^{(9)}$ (the case when $e=0$ ), in which 50 primes satisfy the existence conditions for $e=1$. The smallest primes $p$ such that we failed to construct a row-radical $\left(p, L_{3,5}, 1\right)$-DF are $p=2161$ and $p=8461$.

On the other hand, we can consider a row-radical ( $p, L_{5,3}, 1$ )-DF as a columnradical ( $p, L_{3,5}, 1$ )-DF. Then we can conclude, by Theorem 3.2 .9 , that a columnradical $\left(p, L_{3,5}, 1\right)$-DF exists for any prime $p \equiv 1(\bmod 90)$ with $p>N=$ $L(15,5)>1.1810 \times 10^{12}$ (see Table 3.2 ). However, it is still hard to do a computer search for each prime satisfying $\zeta_{5}+1 \in C^{(9)}$ that less than such a huge $N$.

By using an idea similar to the constructions of $(q, 6,1)$-DFs (see Wilson [120) and ( $q, 7,1$ )-DFs (see Chen, Wei, and Zhu [27), we can efficiently reduce the size of $X_{r, k}$ in Theorem 3.2.6, where a criterion similar to that of $Y_{r, k}$ proposed in Theorem 3.2.13 is necessary. For instance, when $r=6$ and $k=5$, we can take $A_{r}=A_{6}=\{1, \alpha\} \cdot\left\{1, \zeta_{3}, \zeta_{3}^{2}\right\}$ in Theorem 3.2.6, so
that $X_{r, k}=X_{6,5}=\left\{1, \zeta_{3}, \zeta_{3}^{2}\right\} \cdot\left(\{1, \alpha\} \cdot Y_{\frac{r}{2}, k} \cup\left\{\alpha-1, \alpha-\zeta_{3}, \alpha-\zeta_{3}^{2}\right\}\right)$, where $Y_{\frac{r}{2}, k}=Y_{3,5}=\left\{\zeta_{3}-1, \zeta_{5}-1, \zeta_{5}^{2}-1\right\}$ as defined in Theorem 3.2.13. Then we can conclude that, for any prime $p \equiv 1(\bmod 270)$ with $p>2.82 \times 10^{8}>Q(9,4)$ such that each element of $Y_{3,5}$ lies in a distinct coset of $\mathcal{C}^{(9)}$, there exists a row-radical ( $p, L_{5,6}, 1$ )-DF. Actually, we have verified that such an $\alpha$ does exist for each admissible prime $p<2^{26}\left(\approx 6.71 \times 10^{7}\right)$ whose corresponding $Y_{3,5}$ does not have any pair of elements that lies in the same coset of $\mathcal{C}^{(9)}$ except when $p \in\{541,1621,3511,6481\}$.

### 3.3 Row-radical $2 \times k$ grid-block difference families with $k \geq 5$

Next, we propose the concepts of good and bad primes for the construction of $\left(q, L_{2, k}, 1\right)$-DF with $k \geq 5$.
Definition 3.3.1. Let $k \geq 5$ be odd. Let $p \equiv 1\left(\bmod k^{2}\right)$ be a prime. If $\frac{1-\zeta_{k}^{i}}{1-\zeta_{k}^{j}} \notin C^{(k)}$ holds for any $1 \leq j<i \leq \frac{k-1}{2}$, then $p$ is said to be a good prime with respect to $k$, otherwise, a bad prime.

In particular, for $k=5$, if $1+\zeta_{5} \in C^{(5)}$, then $p$ is bad. The following Theorem 3.3.2 characterizes a bad prime with respect to 5 from the aspect of algebraic number theory.
Theorem 3.3.2. For a prime $p \equiv 1(\bmod 25)$, $p$ is bad with respect to 5 if and only if there exists a primary prime $\pi$ in $\mathbb{Z}\left[\zeta_{5}\right]$ such that $\mathrm{N}(\pi)=p$ and $\pi \equiv a$ $(\bmod 5)$, where $\mathrm{N}(\pi)$ denotes the norm of $\pi$ and $a$ is a rational integer with $a \not \equiv 0(\bmod 5)$.

Before coming to the proof of Theorem 3.3.2, we first present a brief review of basic definitions and properties on the cyclotomic field $\mathbb{Q}(\zeta)$, where $\zeta$ denotes a primitive $k$ th root of unity for an odd prime $k$. We make use the set $\left\{\zeta^{i} \mid 1 \leq\right.$ $i \leq k-1\}$ as an integral basis of $\mathbb{Q}(\zeta)$. The conjugate mappings of $\mathbb{Q}(\zeta) / \mathbb{Q}$ are given by $\sigma_{i}(\zeta)=\zeta^{i}$ for $1 \leq i \leq k-1$. Hence, for an element $\xi=a_{1} \zeta+a_{2} \zeta^{2}+$ $\cdots+a_{k-1} \zeta^{k-1}$ with $a_{i} \in \mathbb{Q}(1 \leq i \leq k-1)$, the conjugate can be expressed by

$$
\sigma_{i}(\xi)=a_{1} \zeta^{i}+a_{2} \zeta^{2 i}+\cdots+a_{k-1} \zeta^{(k-1) i}
$$

Denote the complex conjugate of $\xi$ by $\bar{\xi}=\sigma_{k-1}(\xi)$. The norm of $\xi$ is defined by $\mathrm{N}(\xi)=\prod_{i=1}^{k-1} \sigma_{i}(\xi)$.

The properties of primary ideals in $\mathbb{Q}(\zeta)$ and Kummer's reciprocity law play important roles in the proof. First, it is important to introduce the notion of primary elements in $\mathbb{Z}[\zeta]$ for Kummer's reciprocity law.

Definition 3.3.3 (primary cyclotomic integers). Let $\lambda$ denote the prime $1-\zeta \in$ $\mathbb{Z}[\zeta]$. An element $\alpha \in \mathbb{Z}[\zeta]$ is said to be primary if there exists $s \in \mathbb{Z}$ such that the following hold:
$\alpha \not \equiv 0 \quad(\bmod \lambda), \quad \alpha \equiv s \quad\left(\bmod \lambda^{2}\right), \quad$ and $\quad \alpha \bar{\alpha} \equiv s^{2} \quad(\bmod k)$.

Remark. For any $\alpha \in \mathbb{Z}[\zeta]$ satisfying $\alpha \not \equiv 0(\bmod \lambda)$, if $k$ is a regular prime (i.e., the class number of $\mathbb{Q}(\zeta)$ is not divisible by $k$ ), then there exists a unit $u \in \mathbb{Z}[\zeta]^{\times}$such that $\alpha u$ is primary.

Let $\alpha, \pi \in \mathbb{Z}[\zeta]$, where $\pi \neq \lambda$ is a prime. We denote by $\left(\frac{\alpha}{\pi}\right)_{k}$ the $k$ th power residue symbol of $\alpha$ modulo $\pi$. We now state Kummer's reciprocity law.

Theorem 3.3.4 (Kummer's reciprocity law). Let $\pi, \psi \in \mathbb{Z}[\zeta]$ be two distinct primary elements with $(\pi, \psi)=1$. Then,

$$
\left(\frac{\psi}{\pi}\right)_{k}\left(\frac{\pi}{\psi}\right)_{k}^{-1}=1
$$

Definition 3.3.5 (Kummer's quotients of logarithmic derivatives). Let $\alpha=$ $a_{0}+a_{1} \zeta+a_{2} \zeta^{2}+\cdots+a_{k-1} \zeta^{k-1}$ be an element in $\mathbb{Z}[\zeta]$ with $\lambda \nmid \alpha$. Let

$$
\alpha(x)=a_{0}+a_{1} x+a_{2} x^{2}+\cdots+a_{k-1} x^{k-1} \in \mathbb{Z}[x],
$$

then Kummer's quotients of logarithmic derivatives are defined by

$$
\ell_{i}(\alpha)= \begin{cases}\left.\frac{\mathrm{d}^{i} \log \alpha\left(e^{x}\right)}{\mathrm{d} x^{i}}\right|_{x=0} & \text { if } i=1,2, \ldots, k-2 \\ \left.\frac{\mathrm{~d}^{k-1} \log \alpha\left(e^{x}\right)}{\mathrm{d} x^{k-1}}\right|_{x=0} ^{1}+\frac{\alpha(1)-1}{k} & \text { if } i=k-1, \alpha \equiv 1 \quad(\bmod \lambda)\end{cases}
$$

Remark. Kummer's quotients of logarithmic derivatives $\ell_{i}(\alpha)$ are independent of the representation of $\alpha$. In other words, $\ell_{i}(\alpha)$ is uniquely determined by modulo $k$ for each $1 \leq i \leq k-1$.

Theorem 3.3.6 (Kummer's complementary law). If $\alpha \in \mathbb{Z}[\zeta]$ is a primary prime, then $\left(\frac{k}{\alpha}\right)_{k}=\zeta^{\frac{\ell_{k}(\alpha)}{k}}$. (If $\alpha \equiv 1(\bmod \lambda)$, then $\alpha$ is not necessarily primary.) Moreover,

$$
\begin{equation*}
\left(\frac{u}{\alpha}\right)_{k}=\zeta^{\ell_{1}(u) \frac{\mathrm{N}(\alpha)-1}{k}+\sum_{i=1}^{\frac{k-3}{2}} \ell_{2 i}(u) \ell_{k-2 i}(\alpha)} \tag{3.12}
\end{equation*}
$$

for any $u \in \mathbb{Z}[\zeta]^{\times}$.
By employing Kummer's reciprocity law (Theorem 3.3.4) and complementary law (Theorem 3.3.6), we prove Theorem 3.3.2 as follows:

Proof. (Proof of Theorem 3.3.2). Since $\mathbb{Z}[\zeta]$ is a principal ideal domain, there must exist a prime $\pi$ such that $\mathfrak{p}=(\pi)$. Since 5 is a regular prime, without loss of generality, we suppose $\pi$ is a primary prime. Then, $p$ is bad if and only if $\left(\frac{\zeta+1}{\pi}\right)_{5}=1$. Recall that $\mathrm{N}(\pi)=p \equiv 1\left(\bmod 5^{2}\right)$, thus $\zeta^{\frac{\mathrm{N}(\pi)-1}{5}}=1$. By substituting $u=1+\zeta$ and $\alpha=\pi$ in 3.12 of Kummer's complementary law (Theorem 3.3.6), we obtain the following equivalent criterion for $p$ to be a bad prime:

$$
\begin{equation*}
\left(\frac{\zeta+1}{\pi}\right)_{5}=\zeta^{\ell_{2}(1+\zeta) \ell_{3}(\alpha)}=1 \tag{3.13}
\end{equation*}
$$

Now we calculate the exponent explicitly. The first factor is

$$
\begin{equation*}
\ell_{2}(1+\zeta)=\left.\frac{\mathrm{d}^{2} \log \left(1+e^{x}\right)}{\mathrm{d} x^{2}}\right|_{x=0}=\frac{1}{4} \equiv-1 \quad(\bmod 5) \tag{3.14}
\end{equation*}
$$

Next, we suppose $\pi=a_{0}+a_{1} \zeta+a_{2} \zeta^{2}+a_{3} \zeta^{3}$, then the second factor is

$$
\ell_{3}(\pi)=\left.\frac{\mathrm{d}^{3} \log \left(a_{0}+a_{1} e^{x}+a_{2} e^{2 x}+a_{3} e^{3 x}\right)}{\mathrm{d} x^{3}}\right|_{x=0}
$$

On the other hand, $\pi$ is primary if and only if there exist $s, t \in \mathbb{Z}$ with $s \not \equiv 0$ $(\bmod 5)$, such that the following hold:

$$
a_{0} \equiv s+2 t, \quad a_{1} \equiv-t, \quad a_{2} \equiv t, \quad a_{3} \equiv-2 t \quad(\bmod 5)
$$

Therefore,

$$
\begin{equation*}
\ell_{3}(\pi) \equiv \frac{-2 s^{2} t}{s^{3}}=\frac{-2 t}{s} \quad(\bmod 5) \tag{3.15}
\end{equation*}
$$

In summary, it can be derived from (3.13), 3.14, and (3.15) that $\left(\frac{\zeta+1}{\pi}\right)_{5}=$ $\zeta^{\frac{2 t}{s}}=1$, which implies $t \equiv 0(\bmod 5)$. Therefore, there must exist $a \in \mathbb{Z}$ with $a \not \equiv 0(\bmod 5)$ such that $\pi \equiv a(\bmod 5)$ when $p$ is a bad prime with respect to 5.

In general, it is possible to give the explicit condition for a prime $p \equiv 1$ $\left(\bmod k^{2}\right)$ to be a bad prime with respect to any odd prime $k$. However, the calculation would be complicated and the results may not have a simple form of congruences. For instance, when $k=7$, we denote by $\zeta$ the 7 th root of unity. Then, a necessary condition for a prime $p \equiv 1\left(\bmod 7^{2}\right)$ to be a bad prime with respect to 7 is

$$
\pi \equiv\left\{\begin{array}{l}
a+b\left(\zeta+3 \zeta^{3}-2 \zeta^{4}+\zeta^{5}\right)  \tag{3.16}\\
a+b\left(\zeta+3 \zeta^{2}+\zeta^{3}-2 \zeta^{5}\right), \text { or } \\
a+b\left(\zeta-\zeta^{2}-\zeta^{3}+2 \zeta^{4}+2 \zeta^{5}\right)
\end{array} \quad(\bmod 7)\right.
$$

for $a \not \equiv-3 b(\bmod 7)$. For further information on reciprocity laws, algebraic numbers, and cyclotomic fields, the interested reader is referred to 69, [99], [117, and 118.

Example 3.3.7. For $p=1151$, we can choose $\pi=2-5 \zeta_{5}-10 \zeta_{5}^{2}-5 \zeta_{5}^{3}$ as a primary prime in $\mathbb{Z}\left[\zeta_{5}\right]$ which divides 1151 . Since $\pi \equiv 2(\bmod 5)$, it follows from Theorem 3.3.2 that $p=1151$ is a bad prime with respect to 5 .

Example 3.3.8. For $p=1667$, the smallest bad prime with respect to 7 , we can choose $\pi=1-4 \zeta-5 \zeta^{2}-4 \zeta^{3}+7 \zeta^{4}-6 \zeta^{5} \equiv 1-4\left(\zeta+3 \zeta^{2}+\zeta^{3}-2 \zeta^{5}\right)$ $(\bmod 7)$ to meet the criterion in 3.16 .
Remark. All of the bad primes $p \equiv 1(\bmod 25)$ with respect to 5 with $p<10^{4}$ are $1151,1601,1951,3001,3251,3851,4651,4751,5801,6101,7451$, and 9901.

Table 3.3: Parameters for some $\left(p, L_{2,5}, 1\right)$-DDF

| $p$ | $n$ | $e$ | $g$ | $\zeta_{5}$ | $\alpha$ | $p$ | $n$ | $e$ | $g$ | $\zeta_{5}$ | $\alpha$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 101 | 2 | 0 | 2 | 95 | 11 | 3001 | 60 | 1 | 23 | 1125 | 54 |
| 151 | 3 | 0 | 7 | 8 | 13 | 3251 | 65 | 1 | 23 | 1364 | 481 |
| 251 | 5 | 0 | 11 | 219 | 13 | 4751 | 95 | 1 | 19 | 3944 | 346 |

Corollary 3.3.9. There exists a cyclic $\left(p, L_{2,5}, 1\right)$-DDF for every good prime $p \equiv 1(\bmod 50)$.

Proof. For a good prime $p, \zeta_{5}-1$ and $\zeta_{5}^{2}-1$ must lie in different cyclotomic classes, say $\zeta_{5}-1 \in C_{i}^{(5)}$ and $\zeta_{5}^{2}-1 \in C_{j}^{(5)}$, where $i \neq j$ and $i, j \in \mathbb{Z}_{5}$. If there exists $\alpha \in C_{2 i-2 j}^{(5)}$ such that $\alpha-1 \in C_{2 j-1}^{(5)}$, then $X_{2,5}=\left\{\zeta_{5}-1, \zeta_{5}^{2}-\right.$ $\left.1, \alpha\left(\zeta_{5}-1\right), \alpha\left(\zeta_{5}^{2}-1\right), \alpha-1\right\}$ is desired to be a system of representatives of $\mathcal{C}^{(5)}$. This is the case when $e=0$ in Theorem 3.2.6. So it suffices to show the existence of such an element $\alpha \in \mathbb{F}_{p}^{*}$. Then, by setting $n=5$ and $s=2$ in Theorem 1.4 .8 , we have $Q(5,2) \approx 275.6$, which implies a $\left(p, L_{2,5}, 1\right)$-DDF exists for any admissible prime $p>275$. For each of the remaining three primes, namely $p \in\{101,151,251\}$, we list a suitable $\alpha$ in Table 3.3 .

Remark. If we consider the cases when $r=2, k=5$, and $e=1$ in Theorem 3.2.6 then a $\left(p, L_{2,5}, 1\right)$-DDF can be constructed for a bad prime $p \in$ $\{3001,3251,4751\}$. See Table 3.3 for details.

Now we begin to consider a $\left(q, L_{2,5}, 1\right)$-DF when $q$ is a prime power but is not necessarily a prime.

By Theorem 3.2.6, for $r=2$ and $k=5$, it is necessary to check if $1+\zeta_{5}$ is a $5^{e+1}$ th power in $\mathbb{F}_{q}^{*}$. In other words, we have to investigate if $\left(1+\zeta_{5}\right)^{\frac{q-1}{5^{e+1}}}$ is 1 . Let $q=p^{s}$ with prime $p$. The congruence $q \equiv 1(\bmod 50)$ holds exactly in the following cases:
(i) $p \equiv 1(\bmod 50)$ and $s$ arbitrary,
(ii) $p \equiv-1(\bmod 50)$ and $s$ even,
(iii) $p \equiv \pm 7(\bmod 50)$ and $4 \mid s$,
(iv) $p \equiv 1(\bmod 10)$ and $5 \mid s$,
(v) $p \equiv-1(\bmod 10)$ and $10 \mid s$.

The following Lemma 3.3.10 characterizes each case.
Lemma 3.3.10. The following hold when $p$ is a prime:
(i) For $p \equiv 1(\bmod 50)$, Theorem 3.2.6 gives a $\left(p^{s}, L_{2,5}, 1\right)-D F$ if and only if it gives a $\left(p, L_{2,5}, 1\right)-D F$.
(ii) For $p \equiv-1(\bmod 50)$, Theorem 3.2 .6 gives no $\left(p^{2 s}, L_{2,5}, 1\right)-D F$.

Table 3.4: Cyclic ( $p, L_{2, k}, 1$ )-DDFs for $k \geq 7$ and $p<2 \times 10^{4}$

| $k$ | $p$ | $n$ | $g$ | $\zeta_{k}$ | $\alpha$ | $k$ | $p$ | $n$ | $g$ | $\zeta_{k}$ | $\alpha$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 7 | 491 | 5 | 2 | 138 | 25 | 7 | 883 | 9 | 2 | 707 | 20 |
| 7 | 1471 | 15 | 7 | 785 | 7 | 7 | 2549 | 26 | 2 | 2119 | 5 |
| 7 | 5881 | 60 | 31 | 4332 | 10 | 7 | 6469 | 66 | 2 | 1833 | 12 |
| 7 | 6959 | 71 | 7 | 2841 | 56 | 7 | 7253 | 74 | 2 | 3268 | 7 |
| 7 | 7351 | 75 | 7 | 6671 | 93 | 7 | 8429 | 86 | 2 | 6249 | 2 |
| 7 | 8527 | 87 | 5 | 6472 | 7 | 7 | 11369 | 116 | 3 | 5420 | 13 |
| 7 | 16661 | 170 | 11 | 12349 | 21 | 9 | 2593 | 16 | 7 | 251 | 43 |
| 9 | 3889 | 24 | 11 | 923 | 21 | 9 | 15391 | 95 | 17 | 13218 | 18 |
| 9 | 17659 | 109 | 3 | 10277 | 62 | 11 | 727 | 3 | 5 | 662 | 173 |
| 11 | 17183 | 71 | 5 | 14225 | 44 | 13 | 9803 | 29 | 2 | 2774 | 36 |

(iii) For $p \equiv \pm 7(\bmod 50)$, Theorem 3.2 .6 gives no $\left(p^{4 s}, L_{2,5}, 1\right)-D F$.
(iv) For $p \equiv 1(\bmod 10)$, Theorem 3.2 .6 gives a $\left(p^{5 s}, L_{2,5}, 1\right)-D F$ if and only if it gives a $\left(p^{5}, L_{2,5}, 1\right)-D F$.
(v) For $p \equiv-1(\bmod 10)$, Theorem 3.2 .6 gives no $\left(p^{10 s}, L_{2,5}, 1\right)-D F$.

Proof. (i) and (iv) are direct consequences of Lemma 3.2.8. Similarly, it suffices to consider the following cases by Lemma 3.2.8
(ii) For $q=p^{2}$, since $p+1 \equiv 0(\bmod 50)$, we have $n=\frac{p^{2}-1}{50} \equiv 0(\bmod p-1)$. Moreover, since $\operatorname{gcd}(p-1,5)=1$, we have $\frac{p^{2}-1}{5^{e+1}}=\frac{50 n}{5^{e+1}}=5 r(p-1)$ with $r=\frac{2 n}{5^{e}(p-1)}$. Then, $\left(1+\zeta_{5}\right)^{\frac{q-1}{5 e+1}}=\left(1+\zeta_{5}\right)^{5 r(p-1)}=\left(\left(1+\zeta_{5}\right)^{-1}\left(1+\zeta_{5}^{p}\right)\right)^{5 r}=$ $\left(\left(1+\zeta_{5}\right)^{-1}\left(1+\zeta_{5}^{-1}\right)\right)^{5 r}=\left(\zeta_{5}^{-1}\right)^{5 r}=1$. Therefore, in this case, the criteria in Theorem 3.2.6 cannot be satisfied.
(iii) For $q=p^{4}$, since $p^{2}+1 \equiv 0(\bmod 50)$, we have $n=\frac{p^{4}-1}{50} \equiv 0(\bmod p+1)$. Moreover, since $\operatorname{gcd}(p+1,5)=1$, we have $\frac{p^{4}-1}{5^{e+1}}=\frac{50 n}{5^{e+1}}=5 r(p+1)$ with $r=\frac{2 n}{5^{e}(p+1)}$ is even. Then, $z:=\left(1+\zeta_{5}\right)^{\frac{q-1}{5^{+1}}}=\left(1+\zeta_{5}\right)^{5 r(p+1)}=\left(\left(1+\zeta_{5}\right)(1+\right.$ $\left.\left.\zeta_{5}^{p}\right)\right)^{5 r}$. Thus $z=\left(-\zeta_{5}^{4}\right)^{5 r}=1$ if $p \equiv 7(\bmod 50)$, and $z=\left(-\zeta_{5}^{2}\right)^{5 r}=1$ if $p \equiv-7(\bmod 50)$. Therefore, in this case, Theorem 3.2.6 cannot be used.
(v) For $q=p^{10}$, we have $p^{5} \equiv-1(\bmod 50)$. Then one can obtain $(1+$ $\left.\zeta_{5}\right)^{\frac{q-1}{5^{e+1}}}=\left(1+\zeta_{5}\right)^{\frac{\left(p^{5}\right)^{2}-1}{5^{e+1}}}=1$ by proceeding similarly to case (ii), which indicates that Theorem 3.2.6 cannot be used in this case.

Example 3.3.11. More examples of cyclic $\left(p, L_{2, k}, 1\right)$-DDFs for $k \geq 7$ are shown in Table 3.4 .

### 3.4 Kronecker density related to row-radical $2 \times$ $k$ grid-block difference families

In this section, we consider the existence of row-radical $\left(p, L_{2, k}, 1\right)$-DFs with prime $p$ from the viewpoint of Kronecker density.

Let $K$ be a Galois extension of an algebraic number field $F$. For some $\sigma \in \operatorname{Gal}(K / F)$, let $C_{\sigma}$ denote the conjugate class of $\sigma$, i.e., $C_{\sigma}=\left\{\tau \sigma \tau^{-1} \mid \tau \in\right.$ $\operatorname{Gal}(K / F)\}$. Then we define a set $M_{\sigma}$ of prime ideals in $F$ for given $\sigma$ as follows:

$$
M_{\sigma}=\left\{\mathfrak{P} \cap F \mid \mathfrak{P} \text { is a prime ideal in } K \text { such that } \sigma_{\mathfrak{P}} \in C_{\sigma}\right\}
$$

where $\sigma_{\mathfrak{P}}$ is the Frobenius automorphism of $\mathfrak{P}$ over $K$. Now we investigate the Kronecker density of primes with specific properties by using Chebotarëv's Density Theorem (see [99 §25.3) concerning Galois extensions.

Theorem 3.4.1 (Chebotarëv's Density Theorem). The Kronecker density $\delta\left(M_{\sigma}\right)$ of $M_{\sigma}$ is equal to $\frac{\left|C_{\sigma}\right|}{|\operatorname{Gal}(K / F)|}$, i.e.,

$$
\delta\left(M_{\sigma}\right)=\lim _{s \rightarrow 1+0} \sum_{\mathfrak{p} \in M_{\sigma}} \frac{1}{\mathrm{~N}(\mathfrak{p})} / \log \frac{1}{s-1}=\frac{\left|C_{\sigma}\right|}{|\operatorname{Gal}(K / F)|}
$$

where $\mathrm{N}(\mathfrak{p})$ is the norm of the prime ideal $\mathfrak{p}$ in $K$. If the extension $K / F$ is abelian, then there exist infinitely many prime ideals in $F$, say $\mathfrak{p}$, such that $(\mathfrak{p}, K / F)=\sigma$ for each $\sigma \in \operatorname{Gal}(K / F)$, and the density of the set of all those prime ideals is equal to $\frac{1}{[K: F]}$, where $(\mathfrak{p}, K / F)$ is the Artin symbol.

For simplicity, we slightly change the notation in this section. Let $\zeta_{0}$ denote a primitive $k^{2}$-th root of unity and $\zeta=\zeta_{0}^{k}$ in $\mathbb{F}_{p}^{*}$. Then, for $p \equiv 1\left(\bmod k^{2}\right)$, we can see that $\left(\frac{\alpha}{\mathfrak{p}}\right)_{k}=1$ if and only if $\left(\mathfrak{p}, \mathbb{Q}\left(\zeta_{0}, \sqrt[k]{\alpha}\right) / \mathbb{Q}\left(\zeta_{0}\right)\right)=$ id.

### 3.4.1 The Kronecker density of "good" primes

In this subsection, we give a number theoretic discussion of the criteria for $p$ to be a "bad" prime. Let $p \equiv 1\left(\bmod k^{2}\right)$ be a prime. For any odd integer $k$, this congruence is equivalent to $p \equiv 1\left(\bmod 2 k^{2}\right)$.

In Corollary 3.3.9, we settled a necessary condition for the construction of ( $p, L_{2,5}, 1$ )-DFs (see also Definition 3.3.1):

$$
\begin{equation*}
\zeta+1 \notin C^{(5)} \tag{5}
\end{equation*}
$$

Under the assumption that $p \equiv 1(\bmod 25)$, the primes satisfying the condition ( $\mathbf{C}_{5}$ ) are said to be "good" primes, otherwise, "bad" primes. It follows from the following lemma that, among all the primes satisfying $p \equiv 1(\bmod 50)$, the ratios (densities) of "good" primes and "bad" primes are respectively $\frac{4}{5}$ and $\frac{1}{5}$.
Lemma 3.4.2. For $k=5$, the Kronecker densities of "good" primes and "bad" primes are respectively $\frac{1}{25}$ and $\frac{1}{100}$.

Proof. Denote $F=\mathbb{Q}\left(\zeta_{0}\right)$ and $K=F(\sqrt[5]{1+\zeta})$. Let $\mathfrak{p}$ be a prime ideal in $\mathbb{Z}\left[\zeta_{0}\right]$ lying over $(p)$ and let $\sigma_{\mathfrak{p}}:=(\mathfrak{p}, K / F)$ (the Artin symbol). Then $p \equiv 1(\bmod 25)$ is "good" if, and only if, $\sigma_{\mathfrak{p}}(\sqrt[5]{1+\zeta}) \neq \sqrt[5]{1+\zeta}$. Note that $[K: F]=5$ and $[F: \mathbb{Q}]=\varphi(25)=20$. By Chebotarëv's Density Theorem 3.4.1, the Kronecker density of "bad" primes is

$$
\delta_{b}(5)=\frac{1}{5} \cdot \frac{1}{20}=\frac{1}{100}
$$

Accordingly, the Kronecker density of "good" primes is

$$
\delta_{g}(5)=\left(1-\frac{1}{5}\right) \cdot \frac{1}{20}=\frac{1}{25} .
$$

One can check that, among the first 10,000 primes, there are 101 "bad" primes not satisfying condition ( $\mathrm{C}_{5}$. The ratio is extremely close to $1 / 100$.

In general, let $k \geq 5$ be an odd prime. With the notation of Definition 3.3.1, denote $d=\frac{k-1}{2}$, and $\eta_{i j}=\sqrt[k]{\frac{1-\zeta^{i}}{1-\zeta^{j}}}$ for each $1 \leq j<i \leq d$. For $p \equiv 1\left(\bmod k^{2}\right)$, we consider $F:=\mathbb{Q}\left(\zeta_{0}\right)$ and $K:=F\left(\eta_{21}, \eta_{31}, \ldots, \eta_{d 1}\right)$ in what follows.

It is desired that $H_{k} \cup \alpha H_{k} \cup\{\alpha-1\}$ forms a system of representatives of $\mathcal{C}^{(k)}$ for some $\alpha \in \mathbb{F}_{p}^{*}$, where

$$
H_{k}=\left\{\zeta^{i}-1 \mid i \in\{1,2, \ldots, d\}\right\}
$$

Therefore, as a generalization of $\mathrm{C}_{5}$, we have a necessary condition for $k \geq 5$ as follows:

$$
\begin{equation*}
\eta_{i j}^{k} \notin C_{0}^{(k)}, \quad \text { for any } 1 \leq j<i \leq d \tag{k}
\end{equation*}
$$

The primes satisfying condition ( $\mathbf{C}_{\mathrm{k}}$ are said to be "good" primes with respect to $k$, otherwise, "bad" primes.

Lemma 3.4.3. For an odd prime $k$, let $\delta_{g}(k)$ and $\delta_{b}(k)$ denote the Kronecker densities, respectively, of "good" and "bad" primes with respect to $k$. Then,

$$
\begin{aligned}
\delta_{g}(k) & =\frac{(k-1)!}{(k-d)!\cdot k^{d-1}} \cdot \frac{1}{k(k-1)}, \\
\delta_{b}(k) & =\left(1-\frac{(k-1)!}{(k-d)!\cdot k^{d-1}}\right) \frac{1}{k(k-1)} .
\end{aligned}
$$

Proof. Note that $\eta_{i j}=\frac{\eta_{i 1}}{\eta_{j 1}}$ holds for any $i, j$. Moreover, $\left\{\eta_{21}, \eta_{31}, \ldots, \eta_{d 1}\right\}$ is multiplicatively independent when $k$ is prime. (However, for any $d$-tuple of $\eta_{i j}$ 's, it may not be independent.) Hence, $K / F$ is an abelian extension whose Galois group is of type $(\underbrace{k, k, \ldots, k}_{d-1 \text { times }})$. Accordingly, we have

$$
\operatorname{Gal}(K / F) \cong \prod_{i=2}^{d} \operatorname{Gal}\left(F\left(\eta_{i 1}\right) / F\right)
$$

Moreover, $p \equiv 1\left(\bmod k^{2}\right)$ is "good" with respect to $k$, if and only if $\sigma_{\mathfrak{p}}\left(\eta_{i j}\right) \neq$ $\eta_{i j}$ holds for every $1 \leq j<i \leq d$, where $\mathfrak{p}$ is a prime ideal in $\mathbb{Z}\left[\zeta_{0}\right]$ lying over $(p)$, and $\sigma_{\mathfrak{p}}:=(\mathfrak{p}, K / F)$ (the Artin symbol). By applying Chebotarëv's Density Theorem 3.4.1, we can deduce the Kronecker densities of "good" and "bad" primes with respect to $k$ as follows:

$$
\delta_{g}(k)=\frac{\lambda(k)}{k^{d-1}} \cdot \frac{1}{\varphi\left(k^{2}\right)} \quad \text { and } \quad \delta_{b}(k)=\left(1-\frac{\lambda(k)}{k^{d-1}}\right) \frac{1}{\varphi\left(k^{2}\right)},
$$

where $\lambda(k)$ denotes the number of elements of $\operatorname{Gal}(K / F)$ which do not leave $\eta_{i j}$ fixed for all $1 \leq j<i \leq d$, i.e.,

$$
\lambda(k)=\left|\left\{\sigma \in \operatorname{Gal}(K / F) \mid \sigma\left(\eta_{i j}\right) \neq \eta_{i j}, 1 \leq j<i \leq d\right\}\right|
$$

Suppose $\sigma_{i} \in \operatorname{Gal}\left(F\left(\eta_{i 1}\right) / F\right)$ for $2 \leq i \leq d$. First, there are $(k-1)$ ways to choose a $\sigma_{2}$ which does fix $\eta_{21}$. Next, there are $(k-2)$ ways to choose a $\sigma_{3}$ so that $\left\langle\sigma_{2}, \sigma_{3}\right\rangle$ does not fix $\eta_{31}$ and $\eta_{32}$. Analogously, the number of choices of $\sigma_{d}$ is $(k-d+1)$. Thus, we have $\lambda(k)=(k-1)(k-2) \cdots(k-d+1)=\frac{(k-1)!}{(k-d)!}$ Last, by combining with the fact that $[F: \mathbb{Q}]=\varphi\left(k^{2}\right)=k(k-1)$, we complete the proof.

### 3.4.2 The Kronecker density of "arithmetic" primes

Now we begin to study the condition that $H_{k} \cup \alpha H_{k} \cup\{\alpha-1\}$ forms a system of representatives of $\mathcal{C}^{(k)}$ for an element $\alpha \in \mathbb{F}_{p}^{*}$. In other words, if we denote

$$
R_{k}=\left\{\log _{g} x \quad(\bmod k) \mid x \in H_{k}\right\}
$$

then this is equivalent to considering the following condition:

$$
\begin{equation*}
R_{k} \cup\left(R_{k}+a\right) \cup\{b\}=\mathbb{Z}_{k}, \tag{k}
\end{equation*}
$$

where $a=\log _{g} \alpha, b=\log _{g}(1-\alpha)$ and $\log _{g} x$ denotes the discrete logarithm of $x \in \mathbb{F}_{p}^{*}$ to a primitive element $g \in \mathbb{F}_{p}^{*}$.

It is remarkable that $\left|R_{k}\right|=\frac{k-1}{2}$, provided that the equality $\left(\mathrm{R}_{\mathrm{k}}\right.$ holds, namely, $p$ is a "good" prime with respect to $k$.

In order to characterize the condition $\left(\widehat{R_{k}}\right)$, we first introduce some definitions. For an odd $k$, let $S$ be a subset of $\mathbb{Z}_{k}$ with $|S|=\frac{k-1}{2}$. $S$ is said to be arithmetic if there exists $a \in \mathbb{Z}_{k}$ such that $S \cap(S+a)=\emptyset$. Conversely, if $S \cap(S+a)=\emptyset$ holds, then $S$ is said to be a-arithmetic. Then, the only element, say $b$, in $\mathbb{Z}_{k} \backslash(S \cup(S+a))$ can be uniquely determined, which is called the exceptional element of $S$.

Proposition 3.4.4. Given $a, b \in \mathbb{Z}_{k}$, the $a$-arithmetic subset having $b$ as its exceptional element is unique, that is, $\{a+b, 3 a+b, \ldots,(k-2) a+b\}$.
Proof. It is clear that, if $\operatorname{gcd}(a, k)=1$ (which must hold when $k$ is prime), then $\mathbb{Z}_{k}=\{b, a+b, 2 a+b, \ldots,(k-1) a+b\}$ holds for any $b$. Suppose $S$ is an $a$ arithmetic subset having $b$ as its exceptional element. Then, $b \notin S$. Moreover,
$S^{\prime}:=S+a=\mathbb{Z}_{k} \backslash(S \cup\{b\})$ should hold. Hence, $a+b \notin S^{\prime}$, otherwise, $b \in S$. Thus, $a+b \in S$. Thereby, we consequently have $2 a+b \in S^{\prime}, 3 a+b \in S, \ldots$ Finally, we uniquely determine $S=\{a+b, 3 a+b, 5 a+b, \ldots,(k-2) a+b\}$.

More precisely, we can rewrite $R_{k}$ as follows by avoiding the usage of discrete logarithms:

$$
R_{k}=\left\{s \in \mathbb{Z}_{k} \left\lvert\,\left(\frac{1-\zeta^{i}}{\mathfrak{p}}\right)=\zeta^{s}\right., 1 \leq i \leq \frac{k-1}{2}\right\}
$$

where $\mathfrak{p}$ is a prime ideal in $\mathbb{Z}\left[\zeta_{0}\right]$ lying over $(p)$. A prime $p \equiv 1\left(\bmod k^{2}\right)$ is said to be "arithmetic" with respect to $k$ if $p$ is "good" and $R_{k}$ is arithmetic. The property of being arithmetic is independent of the choice of the ideal $\mathfrak{p}$.

Lemma 3.4.5. For an odd prime $k$, let $\delta_{a}(k)$ denote the Kronecker density of "arithmetic" primes with respect to $k$. Then

$$
\delta_{a}(k)=\frac{d \cdot d!}{k^{d-1}} \cdot \frac{1}{k(k-1)}
$$

Proof. Let $S$ be an $a$-arithmetic subset having $b$ as the exceptional element. Then, $-S+(2 b-a)=-\{a+b, 3 a+b, \ldots,(k-2) a+b\}+(2 b-a)=\{b-$ $2 a, b-4 a, \ldots, b-(k-1) a\} \equiv S(\bmod k)$, i.e., $S$ is also a $(-a)$-arithmetic subset having $b-a$ as the exceptional element. In general, for any arithmetic subsets $T \subset \mathbb{Z}_{k}$, there exist exactly two pairs $(u, w),\left(u^{\prime}, w^{\prime}\right) \in \mathbb{Z}_{k}^{2}$ such that $S=u T+w$ and $S=u^{\prime} T+w^{\prime}$ hold, where $u^{\prime}=-u, w^{\prime}=-w+2 b+a$. In other words, each arithmetic subset has exactly two presentations in terms of the pairs $(a, b)$ and $(-a, b-a)$. Therefore, the number of all arithmetic subsets of $\mathbb{Z}_{k}$ is $\frac{k \varphi(k)}{2}=\frac{k(k-1)}{2}$. By considering the orders of elements in $R_{k}$, we have $\frac{k(k-1)}{2} \cdot\left(\frac{k-1}{2}\right)$ ! different ways to determine which cyclotomic class $\left(1-\zeta^{i}\right)$ lies in. Furthermore, since a cyclic permutation acting on $R_{k}$ (regarded as an ordered set) does not change the value of $\eta_{i j}$ for each $1 \leq j<i \leq \frac{k-1}{2}$, by a similar procedure as in the proof of Lemma 3.4.3, we assert that the ratio of the "arithmetic" primes among the "good" primes is $(d \cdot d!) / \frac{(k-1)!}{(k-d)!}$. Combining with the expression of $\delta_{g}(k)$ in Lemma 3.4.3, we have $\delta_{a}(k)=\frac{d \cdot d!}{k^{d-1}} \cdot \frac{1}{k(k-1)}$.

Remark. For $k=5$, a prime $p$ is "arithmetic" if and only if it is "good".
In particular, given a nonzero $\alpha$ satisfying ( $\mathrm{C}_{\mathrm{k}}$, the condition ( $\mathrm{R}_{\mathrm{k}}$ ) can be satisfied for infinitely many primes. Therefore, the following theorem is straightforward from Lemma 3.4.5.

Theorem 3.4.6. There are infinitely many row-radical $L_{2, k}-D F s$ when $k$ is an odd prime.

### 3.5 Recursive constructions

In this section, we present a recursive construction of $\left(v, L_{r, k}, 1\right)$-DDF by using difference matrices. Let $G$ be an additive group. Let $\Gamma=(V, E)$ be a graph with $V=\left\{x_{1}, x_{2}, \ldots, x_{k}\right\}$. A $k \times|G|$ matrix $M$ with entries from $G$ and row vectors $\boldsymbol{r}_{1}, \boldsymbol{r}_{1}, \ldots, \boldsymbol{r}_{k}$ is called a ( $G, \Gamma, 1$ )-difference matrix ( $D M$ for short) if $\boldsymbol{r}_{i}-\boldsymbol{r}_{j}$ covers all elements of $G$ exactly once for every $\left\{x_{i}, x_{j}\right\} \in E$. Moreover, if $\boldsymbol{r}_{i}$ contains no repeated entries for any $1 \leq i \leq k$, then we say $M$ is a $(G, \Gamma, 1)^{*}$ - DM .

Clearly, the existence of a $(G, \Gamma, 1)$-DM (resp., a $(G, \Gamma, 1)^{*}$-DM) implies that of a $\left(G, \Gamma^{\prime}, 1\right)$-DM (resp., a $\left(G, \Gamma^{\prime}, 1\right)^{*}$-DM) for any subgraph $\Gamma^{\prime}$ of $\Gamma$. When $\Gamma$ is the complete graph $K_{k}$, a $(G, \Gamma, 1)$-DM is also known as a $(G, k, 1)$-DM. It is clear that a $(G, k+1,1)$-DM can give a $\left(G, K_{k}, 1\right)^{*}$-DM. In particular, an $\left(\mathbb{F}_{q}, q, 1\right)$-DM (resp., an $\left(\mathbb{F}_{q}, q-1,1\right)^{*}$-DM) can be created by simply taking the multiplicative table of $\mathbb{F}_{q}$ (and removing the row multiplied by 0 ).

Proposition 3.5.1. For every prime power $q$ and every graph $\Gamma$ of order not greater than $q-1$, there exists an $\left(\mathbb{F}_{q}, \Gamma, 1\right)^{*}$-DM.

By using the notion of $(G, \Gamma, 1)^{*}-\mathrm{DM}$, we can obtain the recursive constructions of $(v, \Gamma, 1)$-DDF, which generalize the recursive constructions of $(v, \Gamma, 1)$ DF (see 20, 22, 60, 61), and ( $v, k, 1$ )-DDF (see [24, 49]).

Theorem 3.5.2. If there exist a $\left(G_{1}, \Gamma, 1\right)-D D F, a\left(G_{2}, \Gamma, 1\right)^{*}-D M$, and $a\left(G_{2}, \Gamma, 1\right)$ $D D F$, then there exists a $\left(G_{1} \oplus G_{2}, \Gamma, 1\right)-D D F$. In particular, if $G_{1}=\mathbb{Z}_{u}$ and $G_{2}=\mathbb{Z}_{v}$, then there exists a $\left(\mathbb{Z}_{u v}, \Gamma, 1\right)$-DDF.

Proof. This is very similar to the proof of [20] Theorem 3.2 and 22 Theorem 7.3 due to Buratti and Pasotti on $(v, \Gamma, 1)$-DF. The disjointness of the resultant DDF is guaranteed by the disjointness of the "ingredient DDFs", and the definition of a $\left(G_{2}, \Gamma, 1\right)^{*}$-DM. Actually, the case when $\Gamma$ is not complete are weaker than that of Fuji-Hara, Miao, and Shinohara [49] Theorem 2.2 on complete sets of DDF, and Chang and Ding [24] Proposition 26 on DDF (for the case of cyclic groups). Thus we omit the proof and refer the readers to the above literature.

Corollary 3.5.3. Let $q_{1}, q_{2}, \ldots, q_{s}$ be prime powers and suppose that there exists a $\left(q_{i}, L_{r, k}, 1\right)$-DDF over $\mathbb{F}_{q_{i}}$ for every $1 \leq i \leq s$. Then there exists a $\left(q_{1} q_{2} \cdots q_{s}, L_{r, k}, 1\right)$-DDF over $\bigoplus_{i=1}^{s} \mathbb{F}_{q_{i}}$. In particular, if $q_{1}, q_{2}, \ldots, q_{s}$ are primes, then there exists a cyclic $\left(q_{1} q_{2} \cdots q_{s}, L_{r, k}, 1\right)$-DDF.

Proof. Since the existence of a $\left(q, L_{r, k}, 1\right)$-DF requires $q \equiv 1(\bmod r k(r+k-2))$, there must be $q-1>r k$ for $r+k \geq 4$. By Proposition 3.5.1, an $\left(\mathbb{F}_{q}, L_{r, k}, 1\right)^{*}$ DM always exists. By repeatedly applying Theorem 3.5.2 together with the DM, we can obtain the desired DDF.
Example 3.5.4 (A cyclic $\left(37 \times 19, L_{2,3}, 1\right)$-DF). Define a $\left(\mathbb{Z}_{19}, K_{6}, 1\right)$-DM (which is obviously a $\left(\mathbb{Z}_{19}, L_{2,3}, 1\right)$-DM $)$ by $D=\left(d_{i j}\right)$ with $d_{i j} \equiv i j(\bmod 19)$ for
$1 \leq i \leq 6$ and $1 \leq j \leq 19$ as follows:

$$
\left(\begin{array}{ccccccccccccccccccc}
1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 & 11 & 12 & 13 & 14 & 15 & 16 & 17 & 18 & 0 \\
2 & 4 & 6 & 8 & 10 & 12 & 14 & 16 & 18 & 1 & 3 & 5 & 7 & 9 & 11 & 13 & 15 & 17 & 0 \\
3 & 6 & 9 & 12 & 15 & 18 & 2 & 5 & 8 & 11 & 14 & 17 & 1 & 4 & 7 & 10 & 13 & 16 & 0 \\
4 & 8 & 12 & 16 & 1 & 5 & 9 & 13 & 17 & 2 & 6 & 10 & 14 & 18 & 3 & 7 & 11 & 15 & 0 \\
5 & 10 & 15 & 1 & 6 & 11 & 16 & 2 & 7 & 12 & 17 & 3 & 8 & 13 & 18 & 4 & 9 & 14 & 0 \\
6 & 12 & 18 & 5 & 11 & 17 & 4 & 10 & 16 & 3 & 9 & 15 & 2 & 8 & 14 & 1 & 7 & 13 & 0
\end{array}\right)
$$

Denote $S_{j}=\left[\begin{array}{lll}d_{1 j} & d_{2 j} & d_{3 j} \\ d_{4 j} & d_{5 j} & d_{6 j}\end{array}\right]$ with the entries in the $j$-th column of $D$. By using the base grid-blocks in Example 3.2 .3 , for each $1 \leq j \leq 19$, we define

$$
\mathcal{B}_{j}=\left\{\left[\begin{array}{lll}
1 & 26 & 10 \\
2 & 15 & 20
\end{array}\right]+37 \cdot S_{j}, \quad\left[\begin{array}{ccc}
6 & 8 & 23 \\
12 & 16 & 9
\end{array}\right]+37 \cdot S_{j}\right\}
$$

Next, by using the base grid-block of a cyclic $\left(19, L_{2,3}, 1\right)$-DF, we denote

$$
\mathcal{B}_{0}=\left\{37 \cdot\left[\begin{array}{ccc}
1 & 7 & 11 \\
3 & 2 & 14
\end{array}\right]\right\}
$$

Then, $\bigcup_{j=0}^{19} \mathcal{B}_{j}$ is the set of base grid-blocks of a cyclic $\left(37 \times 19, L_{2,3}, 1\right)$-DF.

## Chapter 4

## Resolvable grid-block coverings

For practical applications, in order to run an experiment such that all the treatments (viz. a parallel class) can be performed simultaneously with each other, resolvability is taken into account. However, "designs" do not always exist, especially when resolvability is desired.

In the case when $r=k$ is odd, Mutoh, Jimbo, and Fu 85] proposed a construction of resolvable $r \times k$ grid-block designs via cyclic $L_{r, k}$-DF with mutually disjoint base grid-blocks. We will generalize their work in Section 4.1.

Moreover, in the application to group testing for identifying all defective items, it is usually desired to test every pair of items at least once. Thus coverings are more desirable than packings. In a resolvable $(v, r \times k, 1)$ gridblock covering, the number of parallel classes $\rho$ should satisfy $\rho \geq\left\lceil\frac{v-1}{r+k-2}\right\rceil$. If the equality holds, the covering is said to be optimal.

In Section 4.2, we will characterize the optimal resolvable $2 \times c$ grid-block coverings for any $c>2$. Then, in Section 4.3, we will prove that an optimal resolvable $2 \times 3$ grid-block covering with $v$ points exists if and only if $v \equiv 0$ $(\bmod 6)$.

### 4.1 Construction of resolvable grid-block designs via grid-block difference families

A $\left(q, L_{r, k}, 1\right)$-DDF has been used for constructions of resolvable grid-block designs and packings by Mutoh et al. 85].
Theorem 4.1.1 (Mutoh et al. 85] Theorem 29). For a prime power $q$, suppose there exists an elementary abelian $\left(q, L_{r, k}, 1\right)-D D F$. If an $(r k, r \times k, 1)$ grid-block design exists, then a resolvable ( $r k q, r \times k, 1$ ) grid-block design exists.

Here Theorem 4.1.1 is generalized to resolvable coverings as well. Furthermore, we allow the "input design", a DDF, to have a non-prime-power order.

This can also be viewed as a generalization of the famous construction for resolvable BIBDs due to Ray-Chaudhuri and Wilson [98] (see also the monograph 50] Theorem 3.2.5).

Theorem 4.1.2. Suppose there exists a cyclic ( $v, L_{r, k}, 1$ )-DDF over (the additive group of) a ring $\mathbf{R}$. Let $\mathbf{R}^{\times}$denote the group of units of $\mathbf{R}$. If
(i) there exist $u_{1}, u_{2}, \ldots, u_{r k} \in \mathbf{R}^{\times}$, such that $u_{i}-u_{j} \in \mathbf{R}^{\times}$for any $1 \leq i<$ $j \leq r k$, and
(ii) there exists an $(r k, r \times k, 1)$ grid-block design (resp., packing, covering),
then a resolvable (rkv, $r \times k, 1$ ) grid-block design (resp., packing, covering) exists.
Proof. Let $M=[r k]$. Suppose $(M, \mathcal{E})$ is an $(r k, r \times k, 1)$ grid-block design (resp., packing, covering). Let $\mathcal{E}=\left\{\mathrm{E}_{0}, \mathrm{E}_{1}, \ldots, \mathrm{E}_{b-1}\right\}$, where $b=\frac{r k-1}{r+k-2}$ (resp., $\left.b \leq \frac{r k-1}{r+k-2}, b \geq \frac{r k-1}{r+k-2}\right)$. Clearly, $\mathrm{E}_{h}$ contains every element in $M$ exactly once for each $0 \leq h \leq b-1$. For any $r k$-set $S=\left\{a_{1}, a_{2}, \ldots, a_{r k}\right\}$, let $\mathrm{E}_{h}(S)$ denote the grid-block obtained by substituting every $i \in[r k]$ with $a_{i} \in S$ in $\mathrm{E}_{h}$.

Let $\mathcal{A}=\left\{\mathrm{A}_{1}, \mathrm{~A}_{2}, \ldots, \mathrm{~A}_{s}\right\}$ be a $\left(v, L_{r, k}, 1\right)$-DDF, where $s=\frac{v-1}{r k(r+k-2)}$. Since $\mathcal{A}$ contains $r k \cdot s=\frac{v-1}{r+k-2}$ distinct elements which are less than $v$, without loss of generality, we can assume that 0 does not appear in any $\mathrm{A}_{\ell}$ for $1 \leq \ell \leq s$.

Now, we begin to construct the resolvable design of order $r k v$ with point-set $V=\mathbf{R} \times M$. We define two types of new grid-blocks,

$$
\begin{aligned}
& \mathrm{C}_{x}^{h}=\mathrm{E}_{h}\left(\left\{\left(x u_{1}, 1\right),\left(x u_{2}, 2\right), \ldots,\left(x u_{r k}, r k\right)\right\}\right), \text { for } x \in \mathbf{R} \text { and } 0 \leq h \leq b-1, \\
& \mathrm{~B}_{\ell}^{j}=\left(u_{j} \cdot \mathrm{~A}_{\ell}\right) \times\{j\}, \text { for } 1 \leq \ell \leq s \text { and } j \in M,
\end{aligned}
$$

whose total number is $v b+r k s=v b+\frac{v-1}{r+k-2}\left(=\frac{r k v-1}{r+k-2}\right.$ if $\mathcal{E}$ forms a design $)$.
For any $g \in \mathbf{R}$, let $\tau_{g}:(a, j) \mapsto(a+g, j)$ be a mapping over $\mathbf{R} \times M$. Let

$$
\begin{aligned}
& \mathcal{C}=\left\{\tau_{g} \mathrm{C}_{x}^{h} \mid g, x \in \mathbf{R}, 0 \leq h \leq b-1\right\} \\
& \mathcal{B}=\left\{\tau_{g} \mathrm{~B}_{\ell}^{j} \mid g \in \mathbf{R}, j \in M, 1 \leq \ell \leq s\right\}
\end{aligned}
$$

Next, we will show that $(V, \mathcal{C} \cup \mathcal{B})$ is an $r \times k$ grid-block design (resp., packing, covering) by calculating the pure (resp., mixed) differences derived from $\mathrm{C}_{x}^{h}$ and $\mathrm{B}_{\ell}^{j}$. For a grid-block B over $V=\mathbf{R} \times M$ and $i, j \in M$, we define a multiset

$$
\partial_{i j} \mathrm{~B}=\{s-t \mid \mathbf{R} \times M \ni(s, i),(t, j): \text { collinear in } \mathrm{B}\},
$$

which is called the pure (resp., mixed) difference list of B when $i=j$ (resp., $i \neq j$ ). For any distinct $i, j \in M$, the pure difference $\partial_{j j} \mathrm{~B}_{\ell}^{i}=\emptyset$. Moreover, $\mathrm{E}_{h}$ contains no duplicated points. So $\partial_{j j} \mathrm{C}_{x}^{h}=\emptyset$ for each $0 \leq h \leq b-1$. On the other hand, since $\mathcal{A}$ is a difference family and $u_{j}$ is a unit in $\mathbf{R}$, we have

$$
\bigcup_{\ell=1}^{s} \partial_{j j} \mathrm{~B}_{\ell}^{j}=u_{j} \bigcup_{\ell=1}^{s} \Delta \mathrm{~A}_{\ell}=u_{j} \Delta \mathcal{A}=u_{j}(\mathbf{R} \backslash\{0\})=\mathbf{R} \backslash\{0\}
$$

Similarly, for any $i, j, k \in M$ with $i \neq j$, it is obvious that $\partial_{i j} \mathrm{~B}_{\ell}^{k}=\emptyset$. On the other hand, for each $0 \leq h \leq b-1$ and each $x \in \mathbf{R}$,

$$
\partial_{i j} \mathrm{C}_{x}^{h}= \begin{cases}\left\{\left(u_{i}-u_{j}\right) x\right\} & \text { if } i, j \text { are collinear in } \mathrm{E}_{h} \\ \emptyset & \text { otherwise }\end{cases}
$$

If $(M, \mathcal{E})$ is a design, by recalling that $u_{i}-u_{j} \in \mathbf{R}^{\times}$, we have

$$
\bigcup_{x \in \mathbb{Z}_{v}} \bigcup_{h=0}^{b-1} \partial_{i j} \mathrm{C}_{x}^{h}=\bigcup_{x \in \mathbb{Z}_{v}}\left\{\left(u_{i}-u_{j}\right) x\right\}=\left(u_{i}-u_{j}\right) \mathbf{R}=\mathbf{R}
$$

Therefore, $(V, \mathcal{C} \cup \mathcal{B})$ is a design. When $(M, \mathcal{E})$ is a packing, if $i, j \in M$ are not collinear in any $\mathrm{E}_{h} \in \mathcal{E}$, then $\bigcup_{h=0}^{b-1} \partial_{i j} \mathrm{C}_{x}^{h}=\emptyset$. In this case, $(V, \mathcal{C} \cup \mathcal{B})$ is a packing. When $(M, \mathcal{E})$ is a covering, if $i, j \in M$ are collinear in $m(m \geq 2)$ gridblocks in $\mathcal{E}$, then $\bigcup_{x \in \mathbb{Z}_{v}} \bigcup_{h=0}^{b-1} \partial_{i j} \mathrm{C}_{x}^{h}$ (as a multiset) consists of all the elements of $\mathbb{Z}_{v}$ with multiplicity $m$. Thus, $(V, \mathcal{C} \cup \mathcal{B})$ is a covering.

It remains to show the resolvability. Let

$$
\begin{align*}
& \mathcal{P}_{g}=\left\{\tau_{g} \mathrm{~B}_{\ell}^{j} \mid j \in M, 1 \leq \ell \leq s\right\} \cup\left\{\tau_{g} \mathrm{C}_{x}^{0} \mid x \notin \bigcup_{i=1}^{s} \mathrm{~A}_{i}\right\}, \text { for each } g \in \mathbf{R},  \tag{4.1}\\
& \mathcal{R}_{x}^{0}=\left\{\tau_{g} \mathrm{C}_{x}^{0} \mid g \in \mathbf{R}\right\}, \text { for each } x \in \bigcup_{i=1}^{s} \mathrm{~A}_{i},  \tag{4.2}\\
& \mathcal{R}_{x}^{h}=\left\{\tau_{g} \mathrm{C}_{x}^{h} \mid g \in \mathbf{R}\right\}, \text { for each } x \in \mathbf{R} \text { and } 1 \leq h \leq b-1 \tag{4.3}
\end{align*}
$$

Then, 4.1), 4.2, and 4.3) give rise to $\rho:=v+s r k+v(b-1)=v b+\frac{v-1}{r+k-2}$ resolution classes. In particular, if $\mathcal{E}$ forms a design, then $\rho=\frac{r k v-1}{r+k-2}$.
Remark. Let $\rho_{-}$(resp., $\rho_{+}$) denote the number of resolution classes of the (vkr, $k \times r, 1$ ) grid-block packing (resp., covering) we constructed. Then, $\rho_{-}=$ $v\left\lfloor\frac{r k-1}{r+k-2}\right\rfloor+\frac{v-1}{r+k-2}$ (resp., $\rho_{+}=v\left\lceil\frac{r k-1}{r+k-2}\right\rceil+\frac{v-1}{r+k-2}$ ). A maximal packing (resp., minimal covering) should satisfy $\rho_{-}=\tilde{\rho}_{-}:=\left\lfloor\frac{r k v-1}{r+k-2}\right\rfloor$ (resp., $\rho_{+}=\tilde{\rho}_{+}:=$ $\left\lceil\frac{r k v-1}{r+k-2}\right\rceil$ ). More precisely, for an $r \times k$ grid-block covering, we have

$$
\varrho_{+}=\lim _{v \rightarrow \infty} \frac{\tilde{\rho}_{+}}{\rho_{+}}=\lim _{v \rightarrow \infty} \frac{\left\lceil\frac{r k v-1}{r+k-2}\right\rceil}{v\left\lceil\frac{r k-1}{r+k-2}\right\rceil+\frac{v-1}{r+k-2}}=\frac{\frac{r k}{r+k-2}}{\frac{1}{r+k-2}+\left\lceil\frac{r k-1}{r+k-2}\right\rceil} \leq 1
$$

where equality holds if and only if $r=k$ is odd (hence it becomes a design). For example, when $r=k+2$, we have $\varrho_{+}=\frac{k^{2}+2 k}{k^{2}+3 k+1}$, which is greater than 0.95 if $k \geq 19$. In addition, when $k$ is a large odd integer, Theorems 3.2.6 and 3.2.13 can also give quite a few cyclic DF, which derive "nearly optimal" coverings via Theorem 4.1.2

In order to simplify criterion (i) in Theorem 4.1.2, we need a simple lemma.

Lemma 4.1.3. Let $s$ and $m>1$ be positive integers. Let $q_{1} \leq q_{2} \leq \cdots \leq q_{s}$ be prime powers. For $\mathbf{R}=\bigoplus_{\ell=1}^{s} \mathbb{F}_{q_{\ell}}$, if $m<q_{1}$, then there exist $\boldsymbol{u}_{1}, \boldsymbol{u}_{2}, \ldots, \boldsymbol{u}_{m} \in$ $\mathbf{R}^{\times}$such that $\boldsymbol{u}_{i}-\boldsymbol{u}_{j} \in \mathbf{R}^{\times}$for any $1 \leq i<j \leq m$.

Proof. This is obvious when $s=1$. Suppose $s \geq 2$. Since $m<q_{1} \leq q_{2} \leq \cdots \leq$ $q_{s}$, we can take $m$ distinct nonzero elements arbitrarily in $\mathbb{F}_{q_{\ell}}$ for each $1 \leq \ell \leq s$, say $u_{1}^{(\ell)}, u_{2}^{(\ell)}, \ldots, u_{m}^{(\ell)}$. Let $\boldsymbol{u}_{i}=\left(u_{i}^{(1)}, u_{i}^{(2)}, \ldots, u_{i}^{(s)}\right)$ for each $1 \leq i \leq m$. It is easy to verify that $\boldsymbol{u}_{1}, \boldsymbol{u}_{2}, \ldots, \boldsymbol{u}_{m}$ are the required units in $\mathbf{R}^{\times}$.

By combining Lemma 4.1.3. Corollary 3.2.10. Corollary 3.5.3, and a $(9,3 \times$ 3,1 ) grid-block design (see [57, 85 ] for the existence), we have the following:

Corollary 4.1.4. There exists a resolvable $(9 v, 3 \times 3,1)$ grid-block design if $v=q_{1} q_{2} \cdots q_{s}$ and $q_{i} \equiv 1(\bmod 36)$ is a prime power for every $1 \leq i \leq s$.

### 4.2 Optimal resolvable grid-block coverings

Let $G$ be a (finite, simple, and undirected) graph and let $V$ be a finite set. Suppose $(V, \mathcal{A})$ is a $G$-covering. The excess graph of $(V, \mathcal{A})$ is the multigraph $(V, E)$, where each edge $\{x, y\}$ occurs in $E$ with multiplicity

$$
\mid\{A \in \mathcal{A} \mid\{x, y\} \text { is an edge in } A\} \mid-1
$$

Clearly, an $r \times c$ grid-block covering is equivalent to an $L_{r, c}$-covering. Now we consider a resolvable $L_{2, c}$-covering and simply write an $L_{2, c}$ - RC or a $\left(v, L_{2, c}, 1\right)$ RC for short.

Lemma 4.2.1. Let $(V, \mathcal{A})$ be a $\left(2 c h, L_{2, c}, 1\right)-R C$. Then $(V, \mathcal{A})$ is optimal if and only if its excess graph forms a 1-factor of $K_{2 c h}$ over $V$.

Proof. Since every grid-block contains $c^{2}$ edges, the number of edges in the excess graph of an optimal $L_{2, c}$-RC is $2 h \cdot h \cdot c^{2}-\binom{2 c h}{2}=h c$. For a fixed vertex $x$, there are $c$ edges adjacent to $x$ in a grid-block. The total degree of $x$ in an optimal $L_{2, c}$-RC is $2 h \cdot c$. Whereas, the degree of $x$ in $K_{2 c h}$ is $2 c h-1$. Since $2 c h-(2 c h-1)=1$, the excess graph in an optimal $L_{2, c}$-RC forms a 1-factor of $K_{2 c h}$. Conversely, if the excess graph of $(V, \mathcal{A})$ forms a 1-factor of $K_{2 c h}$, then the total number of grid-blocks in $\mathcal{A}$ is clearly $2 h^{2}$, i.e., $\rho=2 h$.

It is remarkable that, unlike the leave of an optimal packing, which is known to be a $(c-1)$-factor (see [73] Theorem 2.1), the excess graph of an optimal covering is much smaller and independent with $c$. Especially when $c$ is large, it becomes a stronger condition for a covering to reach optimality. In particular, an optimal covering with the smallest possible order does not always exist.

Lemma 4.2.2. There exists an optimal $\left(2 c, L_{2, c}, 1\right)-R C$ if and only if $c \leq 4$.

Proof. Let $V=\mathbb{Z}_{c} \times \mathbb{Z}_{2}=\left\{x_{i} \mid x \in \mathbb{Z}_{c}, i \in \mathbb{Z}_{2}\right\}$. An optimal ( $2 c, L_{2, c}, 1$ )-RC should consist of two grid-blocks, say $\mathrm{B}_{1}$ and $\mathrm{B}_{2}$. Without loss of generality, let

$$
\mathrm{B}_{1}=\left[\begin{array}{llll}
0_{0} & 1_{0} & \cdots & (c-1)_{0} \\
0_{1} & 1_{1} & \cdots & (c-1)_{1}
\end{array}\right]
$$

Let $M_{1}=\left(\mathbb{Z}_{c} \backslash\{0\}\right) \times\{1\}$. In order to cover the pairs of the form $\left\{0_{0}, \mu_{1}\right\}$ with $\mu_{1} \in M_{1}$, at least $c-2$ points in $M$ should be collinear with $0_{0}$ in $\mathrm{B}_{2}$. In this case, at least a $(c-2)$-clique is contained in the excess graph, which contradicts Lemma 4.2.1 when $c \geq 5$. Therefore, an optimal ( $2 c, L_{2, c}, 1$ )-RC does not exist for $c \geq 5$.

For $c=4$, let $B_{1}=\left[\begin{array}{llll}0_{0} & 1_{0} & 2_{0} & 3_{0} \\ 0_{1} & 1_{1} & 2_{1} & 3_{1}\end{array}\right]$ and $B_{2}=\left[\begin{array}{llll}0_{0} & 1_{0} & 3_{1} & 2_{1} \\ 1_{1} & 0_{1} & 2_{0} & 3_{0}\end{array}\right]$. For $c=3$, let $\mathrm{B}_{1}=\left[\begin{array}{lll}0_{0} & 1_{0} & 2_{0} \\ 0_{1} & 1_{1} & 2_{1}\end{array}\right]$ and $\mathrm{B}_{2}=\left[\begin{array}{lll}0_{0} & 1_{1} & 2_{1} \\ 0_{1} & 2_{0} & 1_{0}\end{array}\right]$. Then $\left(V,\left\{\mathrm{~B}_{1}, \mathrm{~B}_{2}\right\}\right)$ is an optimal $\left(2 c, L_{2, c}, 1\right)$-RC for $c \in\{3,4\}$.

In design theory, RGDDs (resolvable group divisible designs) and frames are commonly used for constructions of new designs. For the standard notions of design theory, the reader is referred to [7, 50]. The "block sizes" of RGDDs and frames can be generalized to graph type. In particular, we need to introduce $L_{2, c}$-RGDDs (resolvable group divisible designs) and $L_{2, c}$-frames (see also [73]).

Let $K_{g_{1}, g_{2}, \ldots, g_{u}}$ denote a complete $u$-partite graph with vertex set $V$ whose partite sets are $G_{1}, G_{2}, \ldots, G_{u}$ with $\left|G_{i}\right|=g_{i}$ for every $1 \leq i \leq u$. Suppose $\mathcal{A}$ is an $L_{2, c}$-decomposition of $K_{g_{1}, g_{2}, \ldots, g_{u}}$. Let $\mathcal{G}=\left\{G_{1}, G_{2}, \ldots, G_{u}\right\}$. Then $(V, \mathcal{G}, \mathcal{A})$ is called an $L_{2, c}$ group divisible design $(G D D)$ of type $\left(g_{1}, g_{2}, \ldots, g_{u}\right)$, where $G_{i}$ is referred to as a group for every $1 \leq i \leq u$. Moreover, if $\mathcal{A}$ can be partitioned into parallel classes, then $(V, \mathcal{G}, \mathcal{A})$ is called a resolvable group divisible design $(R G D D)$. If $g=g_{1}=g_{2}=\cdots=g_{u}$, the GDD is said to be uniform. In this case, $(V, \mathcal{G}, \mathcal{A})$ is called an $L_{2, c}$-GDD of type $g^{u}$ for convenience. Clearly, a (resolvable) $L_{2, c}$-GDD of type $1^{v}$ is nothing but a (resolvable) $\left(v, L_{2, c}, 1\right)$ design.

Proposition 4.2.3 ([73] Lemma 2.2). There are exactly $\frac{g(u-1)}{c}$ parallel classes in any $L_{2, c}$-RGDD of type $g^{u}$.

Let $(V, \mathcal{G}, \mathcal{A})$ be an $L_{2, c^{-}}$GDD with $\mathcal{G}=\left\{G_{1}, G_{2}, \ldots, G_{u}\right\}$. A partial parallel class of $(V, \mathcal{G}, \mathcal{A})$ is a collection of subgraphs of $\mathcal{A}$ whose vertex sets are mutually disjoint. If $\mathcal{A}$ can be partitioned into partial parallel classes, each of which forms a partition of $V \backslash G_{i}$ for some $G_{i} \in \mathcal{G}$, then $(V, \mathcal{G}, \mathcal{A})$ is called an $L_{2, c}$-frame of type $\left(g_{1}, g_{2}, \ldots, g_{u}\right)$, where $\left|G_{i}\right|=g_{i}$ for every $1 \leq i \leq u$. If $g=g_{1}=g_{2}=\cdots=$ $g_{u}$, then $(V, \mathcal{G}, \mathcal{A})$ is called an $L_{2, c}$-frame of type $g^{u}$ for convenience.

Proposition 4.2.4 ([73] Lemma 2.7). Let $(V, \mathcal{G}, \mathcal{A})$ be an $L_{2, c}$-frame. For any $G_{i} \in \mathcal{G}$, the number of partial parallel classes over $V$ missing $G_{i}$ is $\left|G_{i}\right| / c$.

Some fundamental constructions for an optimal $L_{2, c}$-RC by using $L_{2, c}$-RGDDs and $L_{2, c}$-frames are given as follows.

Construction 4.2.5. Let $u$ be an even positive integer. If there exists an $L_{2, c}$-RGDD of type $c^{u}$, then an optimal $\left(c u, L_{2, c}, 1\right)$-RC exists.

Proof. Since $u$ is even, let $(V, \mathcal{G}, \mathcal{A})$ be an $L_{2, c}$-RGDD of type $c^{u}$ with $\mathcal{G}=$ $\left\{G_{1}, G_{2}, \ldots, G_{u / 2}\right\} \cup\left\{H_{1}, H_{2}, \ldots, H_{u / 2}\right\}$. Let $G_{i}=\left\{g_{1}^{(i)}, g_{2}^{(i)}, \ldots, g_{c}^{(i)}\right\}$ and $H_{i}=$ $\left\{h_{1}^{(i)}, h_{2}^{(i)}, \ldots, h_{c}^{(i)}\right\}$ for every $1 \leq i \leq u / 2$. Then, we have $u / 2$ new grid-blocks given by

$$
\mathrm{F}_{i}=\left[\begin{array}{llll}
g_{1}^{(i)} & g_{2}^{(i)} & \ldots & g_{c}^{(i)} \\
h_{1}^{(i)} & h_{2}^{(i)} & \ldots & h_{c}^{(i)}
\end{array}\right] \text { for every } 1 \leq i \leq u / 2
$$

$\mathrm{F}_{i}$ covers all edges in two complete graphs whose vertex sets are $G_{i}$ and $H_{i}$, and

$$
\mathcal{F}_{i}=\left\{\left\{g_{j}^{(i)}, h_{j}^{(i)}\right\} \mid 1 \leq j \leq c\right\}
$$

Then $\left(V, \mathcal{A} \cup\left\{\mathrm{~F}_{1}, \mathrm{~F}_{2}, \ldots, \mathrm{~F}_{u / 2}\right\}\right)$ is an $L_{2, c}-\mathrm{RC}$, in which $\left\{\mathrm{F}_{1}, \mathrm{~F}_{2}, \ldots, \mathrm{~F}_{u / 2}\right\}$ forms one more parallel class. Moreover, $\bigcup_{i=1}^{u / 2} \mathcal{F}_{i}$ is a partition of $V$ into pairs, which is nothing but the excess graph of $\left(V, \mathcal{A} \cup\left\{\mathrm{~F}_{1}, \mathrm{~F}_{2}, \ldots, \mathrm{~F}_{u / 2}\right\}\right)$. By Lemma 4.2.1. $\left(V, \mathcal{A} \cup\left\{\mathrm{~F}_{1}, \mathrm{~F}_{2}, \ldots, \mathrm{~F}_{u / 2}\right\}\right)$ is a desired optimal $L_{2, c}-\mathrm{RC}$.

Construction 4.2.6. Let $g$ be a multiple of $2 c$. Suppose a $\left(g, L_{2, c}, 1\right)$-RC exists. If there exists an $L_{2, c}$-RGDD of type $g^{u}$, then an optimal $\left(g u, L_{2, c}, 1\right)$-RC exists.

Proof. Let $(V, \mathcal{G}, \mathcal{A})$ be an $L_{2, c}$-RGDD of type $g^{u}$ with $\mathcal{G}=\left\{G_{1}, G_{2}, \ldots, G_{u}\right\}$. Suppose $\left(G_{i}, \mathcal{B}_{i}\right)$ is an optimal $\left(g, L_{2, c}, 1\right)$-RC for every $1 \leq i \leq u$. Then every $\left(G_{i}, \mathcal{B}_{i}\right)$ has the same number of parallel classes, because the $\left|G_{i}\right|$ are identical with each other for $1 \leq i \leq u$. Hence $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i}\right)$ is resolvable, where $V=\bigcup_{i=1}^{u} G_{i}$. Clearly, $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i} \cup \mathcal{A}\right)$ is an $L_{2, c}-\mathrm{RC}$ as well. Moreover, by Lemma 4.2.1, the excess graph of $\left(G_{i}, \mathcal{B}_{i}\right)$ forms a 1 -factor of $K_{g}$ on $G_{i}$ for each $1 \leq i \leq u$. The excess graph of $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i} \cup \mathcal{A}\right)$ is obviously the union of 1 -factors, each of which consists of a partition of $G_{i}$. Thus, $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i} \cup \mathcal{A}\right)$ is also optimal by Lemma 4.2.1.

Construction 4.2.7. Let $g$ be a multiple of $2 c$. Suppose a $\left(g, L_{2, c}, 1\right)$-RC exists. If there exists an $L_{2, c}$-frame of type $g^{u}$, then an optimal $\left(g u, L_{2, c}, 1\right)$-RC exists.

Proof. Let $h=\frac{g}{2 c}$. Let $(V, \mathcal{G}, \mathcal{A})$ be an $L_{2, c}$-frame of type $g^{u}$ with $\mathcal{G}=$ $\left\{G_{1}, G_{2}, \ldots, G_{u}\right\}$. By Proposition 4.2.4 the number of partial parallel classes over $V$ missing $G_{i}$ is $2 h$ for every $G_{i} \in \mathcal{G}$. On the other hand, it is shown in the proof of Lemma 4.2.1 that the number of parallel classes in any $\left(g, L_{2, c}, 1\right)$-RC is also $2 h$. For $1 \leq i \leq u$, let $\left(G_{i}, \mathcal{B}_{i}\right)$ be a $\left(g, L_{2, c}, 1\right)$-RC. Clearly, $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i} \cup \mathcal{A}\right)$ is an $L_{2, c}$-covering. By combining a parallel class in $\left(G_{i}, \mathcal{B}_{i}\right)$ with a partial parallel class of $(V, \mathcal{G}, \mathcal{A})$ missing $G_{i}$, a new parallel class over $V$ can be obtained. Proceeding similarly for each $1 \leq i \leq u$, we obtain $2 h$ parallel classes over $V$. Therefore, $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i} \cup \mathcal{A}\right)$ is resolvable. Finally, similarly to the proof of Construction 4.2 .6 that the excess graph of $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i} \cup \mathcal{A}\right)$ forms a 1-factor of $K_{g u}$ on $V$. By Lemma 4.2.1, $\left(V, \bigcup_{i=1}^{u} \mathcal{B}_{i} \cup \mathcal{A}\right)$ is an optimal $\left(g u, L_{2, c}, 1\right)$-RC.

Li and Yin [73] described several recursive constructions of $L_{2, c}$-RGDDs and $L_{2, c}$-frames, which generalize the classical RGDDs, frames, and RBIBDs. We will use the following two theorems from 73.

Theorem 4.2.8 ([73] Construction 2.9). Suppose there exists an $L_{2, c}$-frame of type $(\mathrm{mg})^{h}$. If an $L_{2, c}-R G D D$ of type $g^{m+1}$ exists, then an $L_{2, c}-R G D D$ of type $g^{1+m h}$ exists.

Theorem 4.2.9 (73] Construction 2.13). Suppose there exists a $k$-frame of type $g^{u}$. If an $L_{2, c}-R G D D$ of type $m^{k}$ exists, then an $L_{2, c}$-frame of type $(\mathrm{mg})^{u}$ exists.

### 4.3 Optimal resolvable $2 \times 3$ grid-block coverings

Lemma 4.3.1. There exists an optimal $\left(12, L_{2,3}, 1\right)-R C$.
Proof. The grid-blocks of an optimal $\left(12, L_{2,3}, 1\right)$-RC over $\mathbb{Z}_{12}$ are shown as follows, consisting of 4 parallel classes.

$$
\begin{array}{lll}
{\left[\begin{array}{ccc}
0 & 9 & 2 \\
8 & 3 & 4
\end{array}\right],} & {\left[\begin{array}{ccc}
10 & 6 & 11 \\
1 & 7 & 5
\end{array}\right],} & {\left[\begin{array}{ccc}
0 & 6 & 1 \\
11 & 2 & 3
\end{array}\right],\left[\begin{array}{ccc}
9 & 10 & 4 \\
5 & 8 & 7
\end{array}\right],} \\
{\left[\begin{array}{ccc}
0 & 10 & 7 \\
5 & 2 & 3
\end{array}\right],\left[\begin{array}{ccc}
1 & 11 & 4 \\
9 & 8 & 6
\end{array}\right],} & {\left[\begin{array}{ccc}
0 & 3 & 10 \\
4 & 6 & 5
\end{array}\right],\left[\begin{array}{ccc}
8 & 1 & 2 \\
9 & 11 & 7
\end{array}\right] .}
\end{array}
$$

Lemma 4.3.2. There exists an $L_{2,3}-R G D D$ of type $3^{u}$ for $u \in\{8,12,14\}$.
Proof. Let $X_{3 u}=I_{3} \times\left(\mathbb{Z}_{u-1} \cup\{\infty\}\right)$, where $I_{3}=\{1,2,3\}$. For each $u \in$ $\{8,12,14\}$, we define $u / 2$ base grid-blocks, say $\mathrm{B}_{1}, \mathrm{~B}_{2}, \ldots, \mathrm{~B}_{u / 2}$, shown in Table 4.1. Let $\mathcal{G}_{3 u}=\left\{I_{3} \times\{x\} \mid x \in \mathbb{Z}_{u-1} \cup\{\infty\}\right\}$ and $\mathcal{B}_{3 u}=\left\{\mathrm{B}_{i}+(*, j) \mid 1 \leq i \leq\right.$ $\left.u / 2, j \in \mathbb{Z}_{u-1}\right\}$, where the second component is deduced cyclically modulo $u-1$ leaving $\infty$ fixed. Then $\bigcup_{i=1}^{u / 2}\left\{\mathrm{~B}_{i}+(*, j)\right\}$ is a parallel class for every $j \in \mathbb{Z}_{u-1}$. Therefore $\left(X_{3 u}, \mathcal{G}_{3 u}, \mathcal{B}_{3 u}\right)$ is a desired $L_{2,3}$-RGDD of type $3^{u}$.

Lemma 4.3.3. There exists an $L_{2,3}-R G D D$ of type $6^{u}$ for $u \in\{3,5\}$.
Proof. An $L_{2,3}$-RGDD of type $6^{3}$ is given in [73] Lemma 3.6. For $u=5$, take $X_{30}=I_{3} \times\left(\mathbb{Z}_{8} \cup\left\{\infty_{1}, \infty_{2}\right\}\right)$ and $\mathcal{G}_{30}=\left\{I_{3} \times\{x, x+4\} \mid 0 \leq x \leq\right.$ $3\} \cup\left\{I_{3} \times\left\{\infty_{1}, \infty_{2}\right\}\right\}$, where $I_{3}=\{1,2,3\}$. Let

$$
\begin{aligned}
\mathrm{B}_{1} & =\left[\begin{array}{ccc}
(3,0) & (1,5) & \left(1, \infty_{2}\right) \\
(1,2) & (1,3) & (2,0)
\end{array}\right], \mathrm{B}_{2}=\left[\begin{array}{ccc}
(2,1) & (1,0) & (2,2) \\
\left(2, \infty_{1}\right) & (3,1) & (1,7)
\end{array}\right] \\
\mathrm{B}_{3} & =\left[\begin{array}{ccc}
(1,1) & (3,6) & \left(3, \infty_{2}\right) \\
(1,6) & \left(3, \infty_{1}\right) & (2,5)
\end{array}\right], \mathrm{B}_{4}=\left[\begin{array}{lll}
(3,7) & (2,4) & (2,6) \\
(3,5) & (3,2) & (2,3)
\end{array}\right] \\
\mathrm{B}_{5} & =\left[\begin{array}{ccc}
\left(2, \infty_{2}\right) & (2,7) & (3,4) \\
(1,4) & \left(1, \infty_{1}\right) & (3,3)
\end{array}\right],
\end{aligned}
$$

Table 4.1: $L_{2,3}$-RGDD of type $3^{u}$

| $u$ | $\mathrm{~B}_{1}, \mathrm{~B}_{2}, \ldots, \mathrm{~B}_{u / 2}$ |
| :--- | :--- |
| $u=8$ |  |

$$
\begin{aligned}
& \mathrm{B}_{1}=\left[\begin{array}{lll}
(3,0) & (1,2) & (3,4) \\
(2,5) & (2, \infty) & (3,3)
\end{array}\right], \mathrm{B}_{2}=\left[\begin{array}{ccc}
(3,6) & (2,0) & (3, \infty) \\
(3,1) & (2,4) & (1,0)
\end{array}\right] \\
& \mathrm{B}_{3}=\left[\begin{array}{llll}
(1,1) & (3,5) & (1, \infty) \\
(1,4) & (2,2) & (2,3)
\end{array}\right], \mathrm{B}_{4}=\left[\begin{array}{lll}
(1,6) & (2,1) & (3,2) \\
(1,5) & (2,6) & (1,3)
\end{array}\right]
\end{aligned}
$$

$u=12$

$$
\begin{aligned}
& \mathrm{B}_{1}=\left[\begin{array}{lll}
(3,0) & (1,2) & (1,3) \\
(1,8) & (2,7) & (1,6)
\end{array}\right], \mathrm{B}_{2}=\left[\begin{array}{ccc}
(2,8) & (3,1) & (3,9) \\
(2,4) & (3,10) & (3,3)
\end{array}\right] \\
& \mathrm{B}_{3}=\left[\begin{array}{lll}
(1,1) & (3,6) & (2,9) \\
(1,5) & (3,7) & (1,0)
\end{array}\right], \mathrm{B}_{4}=\left[\begin{array}{ccc}
(2,10) & (3,4) & (1, \infty) \\
(2,0) & (2,6) & (1,4)
\end{array}\right], \\
& \mathrm{B}_{5}=\left[\begin{array}{ccc}
(1,9) & (3,8) & (2,1) \\
(3,2) & (1,7) & (2, \infty)
\end{array}\right], \mathrm{B}_{6}=\left[\begin{array}{ccc}
(2,5) & (1,10) & (2,3) \\
(2,2) & (3, \infty) & (3,5)
\end{array}\right] .
\end{aligned}
$$

$u=14$

$$
\begin{aligned}
\mathrm{B}_{1} & =\left[\begin{array}{ccc}
(3,0) & (3,12) & (2,1) \\
(3,7) & (2,3) & (3,4)
\end{array}\right], \mathrm{B}_{2}=\left[\begin{array}{ccc}
(3,3) & (3,11) & (2,6) \\
(1,2) & (2,4) & (2,7)
\end{array}\right] \\
\mathrm{B}_{3} & =\left[\begin{array}{lll}
(2,2) & (2,9) & (1,3) \\
(2,11) & (1,0) & (1,7)
\end{array}\right], \mathrm{B}_{4}=\left[\begin{array}{lll}
(1, \infty) & (3,5) & (2,10) \\
(1,11) & (2, \infty) & (2,12)
\end{array}\right], \\
\mathrm{B}_{5} & =\left[\begin{array}{ccc}
(3,6) & (3,2) & (2,0) \\
(1,10) & (1,8) & (2,5)
\end{array}\right], \mathrm{B}_{6}=\left[\begin{array}{lll}
(1,12) & (1,9) & (3,1) \\
(3,10) & (3,8) & (1,4)
\end{array}\right] \\
\mathrm{B}_{7} & =\left[\begin{array}{lll}
(1,6) & (1,1) & (3,9) \\
(1,5) & (2,8) & (3, \infty)
\end{array}\right] .
\end{aligned}
$$

and $\mathcal{B}_{30}=\left\{\mathrm{B}_{i}+(*, j) \mid 1 \leq i \leq 5, j \in \mathbb{Z}_{8}\right\}$, where the second component is deduced cyclically modulo 8 leaving $\infty_{1}$ and $\infty_{2}$ fixed. Clearly, $\bigcup_{i=1}^{5}\left\{\mathrm{~B}_{i}+(*, j)\right\}$ is a parallel class for every $j \in \mathbb{Z}_{8}$. Then $\left(X_{30}, \mathcal{G}_{30}, \mathcal{B}_{30}\right)$ is a desired $L_{2,3}$-RGDD of type $6^{5}$.

Lemma 4.3.4 ( 73$]$ Lemma 3.7). For any integer $u \geq 4$ and $u \notin\{8,12,14,18\}$, an $L_{2,3}$-frame of type $12^{u}$ exists.

Lemma 4.3.5. An $L_{2,3}$-frame of type $24^{h}$ exists for $h \in\{4,6,7,9\}$.
Proof. The cases when $h \in\{6,7,9\}$ are shown in [73] Lemma 3.8. For $h=4$, take a 3 -frame of type $4^{4}$ (see 34 IV.5.30 for the existence). Then apply Theorem 4.2 .9 with $c=3, k=3, g=4, u=4$, and $m=6$ to the $L_{2,3}$-RGDD of type $6^{3}$ in Lemma 4.3 .3 to complete the proof.

Lemma 4.3.6. There exists an optimal $\left(24, L_{2,3}, 1\right)-R C$.
Proof. By applying Construction 4.2.5 to the $L_{2,3}$-RGDD of type $3^{8}$ in Lemma 4.3.2. we can show the claim.

Theorem 4.3.7. There exists an optimal ( $\left.12 u, L_{2,3}, 1\right)-R C$ for any positive integer $u$.

Proof. This holds for $u=1,2$ by Lemmas 4.3.1 and 4.3.6. By Construction 4.2.7. Lemmas 4.3.1, and 4.3.6, it suffices to find $L_{2,3}$-frames of type $12^{u}$ or $24^{u}$. Lemma 4.3 .4 gives the $L_{2,3}$-frame of type $12^{u}$ for any integer $u \geq 4$ and $u \notin\{8,12,14,18\}$. Lemma 4.3.5 gives the $L_{2,3}$-frame of type $24^{h}$ for $2 h \in\{8,12,14,18\}$. For $u=3$, by applying Construction 4.2 .5 to the $L_{2,3^{-}}$ RGDD of type $3^{12}$ in Lemma 4.3.2, we obtain an optimal ( $36, L_{2,3}, 1$ )-RC.
Theorem 4.3.8. There exists an optimal $\left(12 u+6, L_{2,3}, 1\right)-R C$ for any positive integer $u$.

Proof. By Construction 4.2.6 and Lemma 4.2.2, it suffices to find $L_{2,3}$-RGDDs of type $6^{u}$. An $L_{2,3}$-RGDD of type $6^{u}$ exists for $u \in\{3,5\}$ by Lemma 4.3.3. By applying Theorem 4.2 .8 with $c=3, g=6$, and $m=2$ to the $L_{2,3}$-frames of type $12^{u}$ in Lemma 4.3.4 one can obtain an $L_{2,3}$-RGDD of type $6^{2 u+1}$ for any positive integer $u$ with $u \notin\{2,3,8,12,14,18\}$. Similarly, by applying Theorem 4.2 .8 with $c=3, g=6$, and $m=4$ to the $L_{2,3}$-frames of type $24^{h}$ in Lemma 4.3.5, an $L_{2,3}$-RGDD of type $6^{4 h+1}$ can be obtained for any $h \in\{1,4,6,7,9\}$. For $u=3$, by applying Construction 4.2 .5 to the $L_{2,3}$-RGDD of type $3^{14}$ in Lemma 4.3.2 we obtain the desired RC.

## Chapter 5

## Concluding remarks and further problems

This dissertation is concerned with affine-invariant quadruple systems, gridblock difference families, and resolvable grid-block coverings, which can be considered as special kinds of cyclic 3-designs, generalizations of cyclic 2-designs, and generalizations of resolvable 2-designs, respectively.

In Chapter 2, we developed two series of constructions for affine-invariant quadruple systems, which depend on 1-factors of a graph and a hypergraph, respectively. On one hand, the graph $\operatorname{CG}\left(\Omega_{p}\right)$ plays an essential role for Construction 2.2.6 for $\operatorname{AsSQS}^{A}(2 p)$ in Section 2.2.2 and Construction 2.5.1 for affine-invariant $\operatorname{TQS}(p)$ in Section 2.5 . However, it is still a challenging problem to theoretically show that $\operatorname{CG}\left(\Omega_{p}\right)$ has a 1-factor for any prime $p \equiv 1,5$ $(\bmod 12)$. Since we have clarified the relationship between the graph $\operatorname{CG}\left(\Omega_{p}\right)$ and the group $\operatorname{PSL}(2, p)$, we shall find a new way to challenge the problem via group theory.

Problem 1. Prove the existence of 1-factor of $\operatorname{CG}\left(\Omega_{p}\right)$ for any prime $p \equiv 1$ $(\bmod 4)$ with the help of group theoretic methods.

On the other hand, Constructions 2.2 .20 relies heavily on an edge-colored hypergraph. However, the definition of the hypergraph and its coloring are complicated. Also, less is known about the existence problem of factors (parallel classes) of hypergraphs (designs). Hence, for that special hypergraph, which is also a PBD, we need to make progress towards the following direction:

Problem 2. Characterize the hypergraph (PBD) defined for Constructions 2.2.20 and derive an algebraic or a design-theoretic criterion for the existence of a rainbow 1-factor.

By the recursive constructions presented in Sections 2.3.2 and 2.3.3, we showed for a prime $p \equiv 1,5(\bmod 12)$ that if the criteria developed for Construction 2.2 .6 or Construction 2.2 .20 can be satisfied, then an $\operatorname{AsSQS}\left(2 p^{m}\right)$ exists for any positive integer $m$.

Together with the results obtained by computer search for Construction 2.2.6 (see Corollary 2.2.8), we conclude that an $\operatorname{AsSQS}\left(2 p^{m}\right)$ exists for every prime $p \equiv 1,5(\bmod 12)$ with $p<10^{5}$ and any positive integer $m$. We leave the following as an open problem:

Problem 3. Find an $\operatorname{AsSQS}\left(2 p_{1} p_{2}\right)$ for distinct primes $p_{1}, p_{2} \equiv 1,5(\bmod 12)$, or prove the non-existence of such kind of AsSQSs.

In Section 2.7, new applications of affine-invariant designs are illustrated. The affine-invariant property provides stronger symmetry and so it is expected to have more applications.

In Chapter 3, we proposed an intermediate algebraic consequence for showing the asymptotic existence of "DF-like" designs, and then used it to improve the existence bounds for grid-block difference families over finite fields. However, for many cases, the bounds are still not good enough.

Problem 4. Improve the bound in Theorem 3.1.3 by considering more relaxed restrictions instead of Theorem 3.1.3 (i) and (ii).

In Section 3.4, the concept of Kronecker density for prime numbers are utilized for grid-block difference families. In general, we should consider the following problem:

Problem 5. Consider the Kronecker density with respect to other combinatorial constructions, such as $t$-designs and combinatorial codes.

It is shown in Section 3.2 that there are "bad" primes when the grid-block size is large, and in that case the radical construction does not work. In order to settle the existence for grid-block difference families, we still need to solve the following:

Problem 6. Provide constructions and existence theorems for the grid-block difference families over the finite fields of "bad" prime orders.

In Chapter 4, we constructed resolvable grid-block designs (packings, coverings) via grid-block difference families. Moreover, we considered the recursive constructions for optimal resolvable grid-block coverings by using frames and RGDDs. Indeed, all the RGDDs in Section 4.3 are computed with the aid of a SAT-based constraint solver - Sugar (cf. [112]; see http://bach.istc.kobeu.ac.jp/sugar/). We can observe that all those RGDDs admit a "rotational" type automorphism, so we should consider the combinatorial nature of those designs.

# List of papers related to this dissertation 

- X.-N. Lu, On affine-invariant two-fold quadruple systems, Graphs and Combinatorics, 31(6): 1915-1927, 2015.
- X.-N. Lu, M. Jimbo, Affine-invariant strictly cyclic Steiner quadruple systems, Designs, Codes and Cryptography, in press.
- X.-N. Lu, Optimal resolvable $2 \times c$ grid-block coverings, Utilitas Mathematica, in press.


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