Operational notes

- 2 Document updated on March 27, 2022.
- The following colors are **not** part of the final product, but serve as highlights in the edit-
- 4 ing/review process:
- text that needs attention from the Subject Matter Experts: Mirco, Anna,& Jan
- terms that have not yet been defined in the book
- things that need to be checked only at the very final typesetting stage
- (and it doesn't make sense to do them before)
- text that needs advice from the communications/marketing team: Aaron & Shane
- text that needs to be completed or otherwise edited (by Sylvia)

Todo list

12	zero-knowledge proofs
13	finite field
14	elliptic curve
15	add references
16	methatical
17	We haven't really talked about what a ring is at this point
18	What's the significance of this distinction?
19	reverse
20	Turing machine
21	polynomial time
22	sub-exponentially, with $\mathcal{O}((1+\varepsilon)^n)$ and some $\varepsilon > 0$
23	Add text
24	\mathbb{Q} of fractions
25	pseudocode
26	check reference
27	modular arithmetics
28	actual division
29	multiplicative inverses
30	factional numbers
31	exponentiation function
32	See XXX
33	This Sage snippet should be described in more detail
34	prime fields
35	residue class rings
36	check algorithm floating
37	(-1) should be (-a)?
38	add reference
39	subtrahend
40	minuend
41	what does this mean?
42	Def Subgroup, Fundamental theorem of cyclic groups
43	add reference
44	Add real-life example of 0?
45	add reference
46	check reference
47	check references to previous examples
48	RSA crypto system
49	size 2048-bits
50	check reference

51	add reference: 27?
52	check reference
53	polynomial time
54	exponential time
55	TODO: Fundamental theorem of finite cyclic groups
56	check reference
57	runtime complexity
58	add reference
59	S: what does "efficiently" mean here?
60	computational hardness assumptions
61	check reference
62	check reference
63	explain last sentence more
64	"equation"?
65	check reference
66	what's the difference between \mathbb{F}_p^* and \mathbb{Z}_p^* ?
67	Legendre symbol
68	Euler's formular
69	These are only explained later in the text, '4.27'
70	are these going to be relevant later?
71	TODO: theorem: every factor of order defines a subgroup
72	Is there a term for this property?
73	a few examples?
74	check reference
75	TODO: DOUBLE CHECK THIS REASONING
76	Mirco: We can do better than this
77	check reference
78	add reference
79	pseudorandom
80	oracle
81	check reference
82	add text on this
83	check reference
84	check reference
85	check reference
86	check reference
87	add more examples protocols of SNARK
88	check reference
89	add reference
90	Abelian groups
91	codomain
92	Check change of wording
93	add reference
94	Expand on this?
95	check reference
96	unify \mathbb{Z}_5 and \mathbb{F}_5 across the board?
97	S: are we introducing elliptic curves in section 1 or 2? 61
98	check reference

99	check reference
100	add reference
101	check reference
102	write paragraph on exponentiation
103	add reference
104	check reference
105	add reference
106	group pairings
107	add reference
108	check reference
109	check reference
110	add reference
111	TODO: Elliptic Curve asymmetric cryptography examples. Private key, generator,
112	public key
112	add reference
	maybe remove this sentence?
114	
115	1
116	1
117	self-intersections
118	check reference
119	check reference
120	jubjub
121	check reference
122	affine plane
123	check reference
124	check reference
125	check reference
126	sign
127	more explanation of what the sign is
128	check reference
129	S: I don't follow this at all
130	check reference
131	add explanation of how this shows what we claim
132	should this def. be moved even earlier?
133	chord line
134	tangential
135	tangent line
136	remove Q ?
137	where?
138	check reference
139	check reference
140	check reference
141	check reference
142	check reference
143	check reference
143	check reference
145	check reference
145	add term

147	add term
148	add reference
149	cofactor clearing
150	add reference
151	check reference
152	check reference
153	add reference
154	add reference
155	check reference
156	check reference
	check reference
157	check reference
158	
159	
160	Explain how
161	write example
162	check reference
163	add reference
164	check reference
165	add reference
166	check reference
167	add reference
168	check reference
169	add reference
170	check reference
171	add reference
172	add reference
173	add reference
174	check reference
175	check reference
175	Check if following Alg is floated too far
	add reference
177	add reference
178	
179	
180	is the label in LATEX correct here?
181	check reference
182	check reference
183	check reference
184	check reference
185	check reference
186	check reference
187	check reference
188	check reference
189	check reference
190	check reference
191	add reference
192	check reference
193	check reference
194	check reference

195	check reference	92
196	check reference	92
197	change "tiny-jubjub" to "pen-jubjub" throughout?	93
198	check reference	94
199	check reference	94
200	check reference	95
201	either expand on this or delete it	95
202	•	95
203	check reference	95
204	check reference	95
205	check reference	95
	check reference	95
206	check reference	95
207		96
208		
209	check reference	96
210		97
211	check reference	97
212	add reference	97
213		97
214	This needs to be written (in Algebra)	97
215	add reference	97
216	add reference	97
217	check reference	97
218	towers of curve extensions	98
219	check reference	98
220	check reference	98
221	check reference	98
222	check reference	99
223		99
224		99
225		00
226		00
227		00
		00
228		01
229		01
230		02
231		02
232		
233		03
234		03
235		03
236		03
237		04
238	$oldsymbol{arepsilon}$	04
239		05
240		05
241		05
242	check if the algorithm is floated properly	05

243	check reference
244	again?
245	check reference
246	circuit
247	signature schemes
248	this was called "pen-jubjub"
249	check reference
250	add reference
251	check reference
252	add references
253	add reference
254	reference text to be written in Algebra
255	check reference
256	check reference
257	check reference
258	add reference
259	algebraic closures
260	check reference
261	check reference
262	check reference
263	check reference
264	check reference
265	disambiguate
266	add reference
267	unify terminology
268	check reference
269	actually make this a table?
270	exercise still to be written?
271	add reference
272	check reference
273	check reference
274	add reference
275	check reference
276	check reference
277	check reference
278	add reference
279	check reference
280	check reference
281	check reference
282	what does this mean?
283	write up this part
	add reference
284	check reference
285	cyclotomic polynomial
286	Pholaard-rho attack
287	todo
288	why?
289 290	check reference
ムガリ	

291	check reference	20
292	what does this mean?	20
293	add reference	20
294	add reference	20
295	check reference	20
296	check reference	21
297	add reference	22
298	add exercise	
299	check reference	
300	add reference	
301	add reference	
302	add reference	
303	check reference	
304	check reference	
305	add reference	
	add reference	
306	add reference	
307	check reference	
308	add reference	
309		
310	add reference	
311	finish writing this up	
312	add reference	
313	correct computations	
314	fill in missing parts	
315	add reference	
316	check equation	
317	Chapter 1?	
318	"rigorous"?	
319	"proving"?	
320	Add example	
321	Add more explanation	
322	6	29
323		29
324	add reference	30
325	add reference	30
326	check reference	30
327	check reference	30
328	Are we using w and x interchangeably or is there a difference between them? 1	31
329	check reference	31
330	jubjub	31
331	check reference	31
332	check reference	31
333	check wording	31
334		31
335	check references	32
336		32
337	add reference	32
338		33

339	add reference	<i>i</i> 3
340	check reference	4
341	check reference	4
342	add reference	5
343	add reference	6
344	Schur/Hadamard product	6
345	add reference	6
346	check reference	6
347	check reference	57
348	add reference	8
349	check reference	9
350	check reference	9
351	check reference	
352	check reference	
353	check reference	0
354	add reference	0
355	add reference	
356	check reference	
357	check reference	
358	check reference	
359	check reference	
360	add reference	
361	check reference	
362	add reference	
363	check reference	_
364	check reference	
	check reference	_
365	Should we refer to R1CS satisfiability (p. 139 here?	
366	check reference	
367	add reference	_
368		
369	check reference	
370		-
371		-
372	check reference	
373	check reference	
374	add reference	_
375	"by"?	
376	check reference	
377	check reference	_
378	add reference	
379	add reference	
380	check reference	
381	add reference	
382	clarify language	_
383	check reference	
384	add reference	_
385	check reference	
386	add reference	6

387	check references
388	add references to these languages?
389	check reference
390	check reference
391	check reference
392	check reference
393	check reference
394	check reference
395	check reference
396	check reference
397	check reference
398	add reference
399	check reference
400	add reference
401	add reference
402	check reference
402	check reference
403	check reference
	check reference
405 406	add reference
	check reference
407	check reference
408	"constraints" or "constrained"?
409	check reference
410	"constraints" or "constrained"?
411	add reference
412	"constraints" or "constrained"?
413	
414	add reference
415	check references
416	check reference
417	add reference
418	can we rotate this by 90° ?
419	check reference
420	add reference
421	add reference
422	shift
423	bishift
424	add reference
425	check reference
426	add reference
427	add reference
428	check reference
429	add reference
430	add reference
431	check reference
432	add reference
433	add reference
434	add reference

435	check reference
436	check reference
437	common reference string
438	simulation trapdoor
439	check reference
440	check reference
441	add reference
442	check reference
443	check reference
444	check reference
445	"invariable"?
446	explain why
447	4 examples have the same title. Change it to be distinct
448	check reference
449	add reference
450	check reference
451	add reference
452	add reference
453	add reference
454	check reference
455	add reference
456	add reference
457	check reference
458	check reference
459	add reference
460	add reference
461	check reference
462	add reference
463	add reference
	add reference
464	check reference
465	add reference
466	add reference
467	add reference
468	add reference
469	add reference
470	add reference
471	add reference
472	add reference
473	check reference
474	
475	
476	
477	add reference
478	add reference
479	add reference
480	add reference
481	add reference
100	add reference 200

483	add reference
484	fix error
485	add reference
486	check reference
487	add reference
488	add reference
489	add reference
490	add reference
491	add reference
492	add reference
493	add reference
494	add reference
495	add reference
496	add reference
497	add reference

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TechnoBob and the Least Scruples crew

500 March 27, 2022

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645				The Verification Phase	
645 646 647				Proof Simulation)()

Chapter 1

Introduction

This is dump from other papers as inspiration for the intro:

Zero knowledge proofs are a class of cryptographic protocols in which one can prove honest computation without revealing the inputs to that computation. A simple high-level example of a zero-knowledge proof is the ability to prove one is of legal voting age without revealing the respective age. In a typical zero knowledge proof system, there are two participants: a prover and a verifier. A prover will present a mathematical proof of computation to a verifier to prove honest computation. The verifier will then confirm whether the prover has performed honest computation based on predefined methods. Zero knowledge proofs are of particular interest to public blockchain activities as the verifier can be codified in smart contracts as opposed to trusted parties or third-party intermediaries.

Zero-knowledge proofs (ZKPs) are an important privacy-enhancing tool from cryptography. Theyallow proving the veracity of a statement, related to confidential data, without revealing any in-formation beyond the validity of the statement. ZKPs were initially developed by the academiccommunity in the 1980s, and have seen tremendous improvements since then. They are now ofpractical feasibility in multiple domains of interest to the industry, and to a large community ofdevelopers and researchers. ZKPs can have a positive impact in industries, agencies, and for per-sonal use, by allowing privacy-preserving applications where designated private data can be madeuseful to third parties, despite not being disclosed to them.

ZKP systems involve at least two parties: a prover and a verifier. The goal of the prover is toconvince the verifier that a statement is true, without revealing any additional information. Forexample, suppose the prover holds a birth certificate digitally signed by an authority. In orderto access some service, the prover may have to prove being at least 18 years old, that is, thatthere exists a birth certificate, tied to the identify of the prover and digitally signed by a trusted certification authority, stating a birthdate consistent with the age claim. A ZKP allows this, without the prover having to reveal the birthdate.

1.1 Target audience

This book is accessible for both beginners and experienced developers alike. Concepts are gradually introduced in a logical and steady pace. Nonetheless, the chapters lend themselves rather well to being read in a different order. More experienced developers might get the most benefit by jumping to the chapters that interest them most. If you like to learn by example, then you should go straight to the chapter on Using Clarinet.

It is assumed that you have a basic understanding of programming and the underlying logical concepts. The first chapter covers the general syntax of Clarity but it does not delve into what

programming itself is all about. If this is what you are looking for, then you might have a more difficult time working through this book unless you have an (undiscovered) natural affinity for such topics. Do not let that dissuade you though, find an introductory programming book and press on! The straightforward design of Clarity makes it a great first language to pick up.

1.2 The Zoo of Zero-Knowledge Proofs

```
First, a list of zero-knowledge proof systems:
689
       1. Pinocchio (2013): Paper
690

    Notes: trusted setup

691
       2. BCGTV (2013): Paper
692
             - Notes: trusted setup, implementation
693
       3. BCTV (2013): Paper
694
             - Notes: trusted setup, implementation
695
       4. Groth16 (2016): Paper
696

    Notes: trusted setup

697
             - Other resources: Talk in 2019 by Georgios Konstantopoulos
698
       5. GM17 (207): Paper
699

    Notes: trusted setup

700
             - Other resources: later Simulation extractability in ROM, 2018
701
       6. Bulletproofs (2017): Paper
702

    Notes: no trusted setup

703
             - Other resources: Polynomial Commitment Scheme on DL, 2016 and KZG10, Poly-
704
                nomial Commitment Scheme on Pairings, 2010
705
       7. Ligero (2017): Paper
706
             - Notes: no trusted setup
707
             - Other resources:
708
       8. Hyrax (2017): Paper
709

    Notes: no trusted setup

710
             – Other resources:
711
       9. STARKs (2018): Paper
```

- Notes: no trusted setup

– Other resources:

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```
10. Aurora (2018): Paper
715

    Notes: transparent SNARK

716
             – Other resources:
717
      11. Sonic (2019): Paper
718
             - Notes: SNORK - SNARK with universal and updateable trusted setup, PCS-based
719
             - Other resources: Blog post by Mary Maller from 2019 and work on updateable and
720
                universal setup from 2018
721
      12. Libra (2019): Paper
722

    Notes: trusted setup

723
             - Other resources:
      13. Spartan (2019): Paper
725

    Notes: transparent SNARK

726
             - Other resources:
727
      14. PLONK (2019): Paper
728

    Notes: SNORK, PCS-based

729

    Other resources: Discussion on Plonk systems and Awesome Plonk list

730
      15. Halo (2019): Paper
731

    Notes: no trusted setup, PCS-based, recursive

732
             – Other resources:
733
      16. Marlin (2019): Paper
734

    Notes: SNORK, PCS-based

735

    Other resources: Rust Github

736
      17. Fractal (2019): Paper
737

    Notes: Recursive, transparent SNARK

738
             – Other resources:
739
      18. SuperSonic (2019): Paper
740

    Notes: transparent SNARK, PCS-based

741

    Other resources: Attack on DARK compiler in 2021

742
      19. Redshift (2019): Paper
743
             - Notes: SNORK, PCS-based
744
             – Other resources:
745
```

Other resources on the zoo:

reference document

746

Awesome ZKP list on Github, ZKP community with the

748 To Do List

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- Make table for prover time, verifier time, and proof size
- Think of categories Achieved Goals: Trusted setup or not, Post-quantum or not, ...
- Think of categories *Mathematical background*: Polynomial commitment scheme, ...
- ... while we discuss the points above, we should also discuss a common notation/language for all these things. (E.g. transparent SNARK/no trusted setup/STARK)

754 Points to cover while writing

- Make a historical overview over the "discovery" of the different ZKP systems
- Make reader understand what paper is build on what result etc. the tree of publications!
- Make reader understand the different terminology, e.g. SNARK/SNORK/STARK, PCS, R1CS, updateable, universal, ...
- Make reader understand the mathematical assumptions and what this means for the zoo.
 - Where will the development/evolution go? What are bottlenecks?

761 Other topics I fell into while compiling this list

- Vector commitments: https://eprint.iacr.org/2020/527.pdf
- Snarkl: http://ace.cs.ohio.edu/~gstewart/papers/snaarkl.pdf
- Virgo?: https://people.eecs.berkeley.edu/~kubitron/courses/cs262a-F19/projects/reports/project5_report_ver2.pdf

Chapter 2

Preliminaries

2.1 Preface and Acknowledgements

This book began as a set of lecture and notes accompanying the zk-Summit 0x and 0xx It arose from the desire to collect the scattered information of snarks [] and present them to an audience that does not have a strong backgroud in cryptography []

2.2 Purpose of the book

The first version of this book is written by security auditors at Least Authority where we audited quite a few snark based systems. Its included "what we have learned" destilate of the time we spend on various audits.

We intend to let illus- trative examples drive the discussion and present the key concepts of pairing computation with as little machinery as possible. For those that are fresh to pairing-based cryptography, it is our hope that this chapter might be particularly useful as a first read and prelude to more complete or advanced expositions (e.g. the related chapters in [Gal12]).

On the other hand, we also hope our beginner-friendly intentions do not leave any sophisticated readers dissatisfied by a lack of formality or generality, so in cases where our discussion does sacrifice completeness, we will at least endeavour to point to where a more thorough exposition can be found.

One advantage of writing a survey on pairing computation in 2012 is that, after more than a decade of intense and fast-paced research by mathematicians and cryptographers around the globe, the field is now racing towards full matu- rity. Therefore, an understanding of this text will equip the reader with most of what they need to know in order to tackle any of the vast literature in this remarkable field, at least for a while yet.

Since we are aiming the discussion at active readers, we have matched every example with a corresponding snippet of (hyperlinked) Magma [BCP97] code 1, where we take inspiration from the helpful Magma pairing tutorial by Dominguez Perez et al. [DKS09].

Early in the book we will develop examples that we then later extend with most of the things we learn in each chapter. This way we incrementally build a few real world snarks but over full fledged cryptographic systems that are nevertheless simple enough to be computed by pen and paper to illustrate all steps in grwat detail.

2.3 How to read this book

797 Books and papers to read: XXXXXXXXXX

Correctly prescribing the best reading route for a beginner naturally requires individual diagnosis that depends on their prior knowledge and technical preparation.

2.4 Cryptological Systems

The science of information security is referred to as *cryptology*. In the broadest sense, it deals with encryption and decryption processes, with digital signatures, identification protocols, cryptographic hash functions, secrets sharing, electronic voting procedures and electronic money.

EXPAND

6 2.5 SNARKS

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2.6 complexity theory

Before we deal with the mathematics behind zero knowledge proof systems, we must first clarify what is meant by the runtime of an algorithm or the time complexity of an entire mathematical problem. This is particularly important for us when we analyze the various snark systems...

For the reader who is interested in complexity theory, we recommend, or example or, as well as the references contained therein.

813 2.6.1 Runtime complexity

The runtime complexity of an algorithm describes, roughly speaking, the amount of elementary computation steps that this algorithm requires in order to solve a problem, depending on the size of the input data.

Of course, the exact amount of arithmetic operations required depends on many factors such as the implementation, the operating system used, the CPU and many more. However, such accuracy is seldom required and is mostly meaningful to consider only the asymptotic computational effort.

In computer science, the runtime of an algorithm is therefore not specified in individual calculation steps, but instead looks for an upper limit which approximates the runtime as soon as the input quantity becomes very large. This can be done using the so-called *Landau notation* (also called big -\$\mathcal{O}\$-notation) A precise definition would, however, go beyond the scope of this work and we therefore refer the reader to .

For us, only a rough understanding of transit times is important in order to be able to talk about the security of crypographic systems. For example, $\mathcal{O}(n)$ means that the running time of the algorithm to be considered is linearly dependent on the size of the input set n, $\mathcal{O}(n^k)$ means that the running time is polynomial and $\mathcal{O}(2^n)$ stands for an exponential running time (chapter 2.4).

An algorithm which has a running time that is greater than a polynomial is often simply referred to as *slow*.

A generalization of the runtime complexity of an algorithm is the so-called *time complexity* of a mathematical problem, which is defined as the runtime of the fastest possible algorithm that can still solve this problem (chapter 3.1).

Since the time complexity of a mathematical problem is concerned with the runtime analysis of all possible (and thus possibly still undiscovered) algorithms, this is often a very difficult and deep-seated question .

For us, the time complexity of the so-called discrete logarithm problem will be important. This is a problem for which we only know slow algorithms on classical computers at the moment, but for which at the same time we cannot rule out that faster algorithms also exist.

STUFF ON CRYPTOGRAPHIC HASH FUNCTIOND

2.7 Software Used in This Book

2.7.1 Sagemath

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It order to provide an interactive learning experience, and to allow getting hands-on with the concepts described in this book, we give examples for how to program them in the Sage programming language. Sage is a dialect of the learning-friendly programming language Python, which was extended and optimized for computing with, in and over algebraic objects. Therefore, we recommend installing Sage before diving into the following chapters.

The installation steps for various system configurations are described on the sage websit ¹. Note however that we use Sage version 9, so if you are using Linux and your package manager only contains version 8, you may need to choose a different installation path, such as using prebuilt binaries.

We recommend the interested reader, who is not familiar with sagemath to read on the many tutorial before starting this book. For example

https://doc.sagemath.org/html/en/installation/index.html

General Chapter 3

Arithmetics

3.1 Introduction

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3.1.1 Aims and target audience

The goal of this chapter is to enable a reader who is starting out with nothing more than basic high school algebra to be able to solve basic tasks in elliptic curve cryptography without the need of a computer.

How much mathematics do you need to understand zero-knowledge proofs? The answer, of course, depends on the level of understanding you aim for. It is possible to describe zero-knowledge proofs without using mathematics at all; however, to read a foundational paper like Groth [2016], some knowledge of mathematics is needed to be able to follow the discussion.

zeroknowledge proofs

Without a solid grounding in mathematics, someone who is interested in learning the concepts of zero-knowledge proofs, but who has never seen or done computations on, say, a finite field, or an elliptic curve, may quickly become overwhelmed. This is not so much due to the complexity of the mathematics needed, rather because of the vast amount of technical jargon, unknown terms, and obscure symbols that quickly makes a text unreadable, even though the concepts themselves are not actually that hard. As a result, the reader might either lose interest, or pick up some incoherent bits and pieces of knowledge that, in the worst case scenario, result in immature code.

finite field

elliptic curve

This is why we dedicated this chapter to explaining the mathematical foundations needed to understand the basic concepts underlying snark development. We encourage the reader who is not familiar with basic number theory and elliptic curves to take the time and read this and the following chapters, until they are able to solve at least a few exercises in each chapter.

If, on the other hand, you are already skilled in elliptic curve cryptography, feel free to skip this chapter and only come back to it for reference and comparison. Maybe the most interesting parts are XXX .

add references

We start our explanations at a very basic level, and only assume pre-existing knowledge of fundamental concepts like integer arithmetics. At the same time, we'll attempt to teach you to "think mathematically", and to show you that there are numbers and methatical structures out there that appear to be very different from the things you learned about in high school, but on a deeper level, they are actually quite similar.

methatical

We want to stress, however that this introduction is informal, incomplete and optimized to enable the reader to understand zero-knowledge concepts as efficiently as possible. Our focus and design choices are to include as little theory as necessary, focusing on the wealth of numerical examples. We believe that such an informal, example-driven approach to learning

mathematics may make it easier for beginners to digest the material in the initial stages.

For instance, as a beginner, you would probably find it more beneficial to first compute a simple toy **snark** with pen and paper all the way through, before actually developing real-world production-ready systems. In addition, it's useful to have a few simple examples in your head before getting started with reading actual academic papers.

However, in order to be able to derive these toy examples, some mathematical groundwork is needed. This chapter therefore will help you focus on what is important, accompanied by exercises that you are encouraged to recompute yourself. Every section usually ends with a list of additional exercises in increasing order of difficulty, to help the reader memorize and apply the concepts.

3.1.2 The structure of this chapter

We start with a brief recapitulation of basic integer arithmetics like long division, the greatest common divisor and Euclid's algorithm. After that, we introduce modular arithmetics as **the most important skill** to compute our pen-and-paper examples. We then introduce polynomials, compute their analogs to integer arithmetics and introduce the important concept of Lagrange interpolation.

After this practical warm up, we introduce some basic algebraic terms like groups and fields, because those terms are used very frequently in academic papers relating to zero-knowledge proofs. The beginner is advised to memorize those terms and think about them. We define these terms in the general abstract way of mathematics, hoping that the non mathematical trained reader will gradually learn to become comfortable with this style. We then give basic examples and do basic computations with these examples to get familiar with the concepts.

3.2 Integer Arithmetics

In a sense, integer arithmetics is at the heart of large parts of modern cryptography, because it provides the most basic tools for doing computations in those systems. Fortunately, most readers will probably remember integer arithmetics from school. It is, however, important that you can confidently apply those concepts to understand and execute computations in the many pen-and-paper examples that form an integral part of the MoonMath Manual. We will therefore recapitulate basic arithmetics concepts to refresh your memory and fill any knowledge gaps.

In what follows, we apply standard mathematical notations, and use the symbol \mathbb{Z} for the set of all **integers**: S: I think it'd be useful to explain the difference between := and = as well. We have a table on this in the ZKAPs whitepaper.

M:Yeah maybe we use the more suggestive leftarrows aks \leftarrow ? If a table of symbols is unavoidable then ok, I find I super ugly, though

S: The table below is similar to what we have in the ZKAPs whitepaper. I think it's easier to read than a wall of text explaining the same things. We can make it more visually different (e.g. typesetting it in "dark mode", and probably move it earlier in the chapter.

Symbol	Meaning of Symbol	Example	Explanation
=	equals	a=r	a and r have the same value
:=	defining a variable	$\mathbb{Z} := \{\dots, -2, -1, 0, 1, 2, \dots\}$	\mathbb{Z} is the set of integers
\in	from the set	$a \in \mathbb{Z}$	a is an integer

Notation used in this chapter

$$\mathbb{Z} := \{\dots, -3, -2, -1, 0, 1, 2, 3, \dots\} \tag{3.1}$$

Integers are also known as **whole numbers**, that is, numbers that can be written without fractional parts. Examples of numbers that are **not** integers are $\frac{2}{3}$, 1.2 and -1280.006.

If $a \in \mathbb{Z}$ is an integer, then |a| stands for the **absolute value** of a, that is, the the non-negative value of a without regard to its sign:

$$|4| = 4 \tag{3.2}$$

$$|-4| = 4 \tag{3.3}$$

In addition, we use the symbol \mathbb{N} for the set of all **counting numbers** (also called natural numbers). So whenever you see the symbol \mathbb{N} , think of the set of all non negative integers including the number 0:

$$\mathbb{N} := \{0, 1, 2, 3, \dots\} \tag{3.4}$$

We

Any number that is smaller than 0, that is, any number that has a minus sign, is not part of \mathbb{N} ". All counting numbers are integers, but not the other way round. In other words, counting numbers are a subset of integers.

To make it easier to memorize new concepts and symbols, we might frequently link to definitions (See 3.1 for a definition of \mathbb{Z}) in the beginning, but as to many links render a text unreadable, we will assume the reader will become familiar with definitions as the text proceeds at which point we will not link them anymore.

Both sets $\mathbb N$ and $\mathbb Z$ have a notion of addition and multiplication defined on them. Most of us are probably able to do many integer computations in our head, but this gets more and more difficult as these increase in complexity. We will frequently invoke the SageMath system (2.7.1) for more complicated computations. One way to invoke the integer type in Sage is: We haven't

really talked about what a ring is at this point

```
haven't
    sage: ZZ # A sage notation for the integer type
948
    Integer Ring
949
                                                                               talked
    sage: NN # A sage notation for the counting number type
950
                                                                               about
    Non negative integer semiring
951
                                                                               what a
    sage: ZZ(5) # Get an element from the Ring of integers
952
                                                                               ring is at
    5
                                                                              6 this point
953
    sage: ZZ(5) + ZZ(3)
954
                                                                              8
955
    sage: ZZ(5) * NN(3)
                                                                              9
956
    15
                                                                              10
957
    sage: ZZ.random_element(10**50)
                                                                              11
958
    84171518381395594090483054147323086275118949345354
                                                                              12
959
```

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```
sage: ZZ(27713).str(2) # Binary string representation
                                                                           13
960
   110110001000001
                                                                           14
961
    sage: NN(27713).str(2) # Binary string representation
                                                                           15
962
    110110001000001
                                                                           16
963
    sage: ZZ(27713).str(16) # Hexadecimal string representation
                                                                           17
964
    6c41
                                                                           18
965
```

One set of numbers that is of particular interest to us is **prime numbers**, which are counting numbers $p \in \mathbb{N}$ with $p \ge 2$, which are only divisible by themself and by 1. All prime numbers apart from the number 2 are called **odd** (since even numbers greater than 2 are all divisible by 2, they are not prime numbers). We write \mathbb{P} for the set of all prime numbers and $\mathbb{P}_{>3}$ for the set of all odd prime numbers. \mathbb{P} is infinite and can be ordered according to size, so that we can write them as follows:

$$2,3,5,7,11,13,17,19,23,29,31,37,41,43,47,53,59,61,67,...$$
 (3.5)

This is sequence A000040 in OEIS, the Online Encyclopedia of Integer Sequences. In particular, we can talk about small and large prime numbers.

As the **fundamental theorem of arithmetics** tells us, prime numbers are, in a certain sense, the basic building blocks from which all other natural numbers are composed. To see that, let $n \in \mathbb{N}_{\geq 2}$ be any natural number. Then there are always prime numbers $p_1, p_2, \dots, p_k \in \mathbb{P}$, such that

What's the significance of this distinction?

$$n = p_1 \cdot p_2 \cdot \ldots \cdot p_k \,. \tag{3.6}$$

This representation is unique for each natural number (except for the order of the factors) and is called the **prime factorization** of n. 979

Example 1 (Prime Factorization). To see what we mean by prime factorization of a number, let's look at the number $19214758032624000 \in \mathbb{N}$. To get its prime factors, we can successively divide it by all prime numbers in ascending order starting with 2:

We can double check our findings invoking Sage, which provides an algorithm to factor count-980 ing numbers: 981

```
sage: n = NN(19214758032624000)
                                                                            19
982
   sage: factor(n)
                                                                           20
983
   2^7 * 3^3 * 5^3 * 7 * 11 * 17^2 * 23 * 43^2 * 47
                                                                           21
984
```

This computation reveals an important observation: Computing the factorization of an integer is computationally expensive, while the reverse, that is, computing the product of given are reverse set of prime numbers, is fast. MIRCO: inverse process? Its not reversing something actually, but my english doesn't resolve to this level very well

From this, an important question arises: How fast we can compute the prime factorization of a natural number? This is the famous factorization problem and, as far as we know, there is no method on a classical Turing machine that is able to compute this representation in polynomial time. The fastest algorithm known today run sub-exponentially, with $\mathcal{O}((1+\varepsilon)^n)$ and some

It follows that number factorization \Leftrightarrow prime number multiplication is an example of a so-called **one-way function**: Something that is easy to compute in one direction, but hard to

Turing

time

machine

polynomial

 ε)ⁿ) and

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998

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1000

1001

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1003

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1009

1010

1013

1014

1015

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1017

1019

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compute in the other direction. The existence of one-way functions is a basic cryptographic assumptions that the security of many crypto systems is based on.

It should be pointed out, however, hat the American mathematician Peter Williston Shor developed an algorithm in 1994 which can calculate the prime factor representation of a natural number in polynomial time on a quantum computer. The consequence of this is that cryptosystems, which are based on the time complexity of the prime factor problem, are unsafe as soon as practically usable quantum computers become available. Add text_along the lines of "this is the best we got for now" Possibly something on when we can reasonably expect quantum computers to become accessible/usable enough

Add text

- Exercise 1. What is the absolute value of the integers -123, 27 and 0?
- Exercise 2. Compute the factorization of 6469693230 and double check your results using Sage.
- Exercise 3. Consider the following equation $4 \cdot x + 21 = 5$. Compute the set of all solutions for x under the following alternative assumptions:
 - 1. The equation is defined over the type of natural numbers.
 - 2. The equation is defined over the type of integers.
- Exercise 4. Consider the following equation $2x^3 x^2 2x = -1$. Compute the set of all solutions x under the following assumptions:
 - 1. The equation is defined over the type of natural numbers.
 - 2. The equation is defined over the type of integers.
 - 3. The equation is defined over the type \mathbb{Q} of fractions.

Q of fractions

Euclidean Division Division in the commonly understood sense is not defined for integers, as, for example, 7 divided by 3 will not be an integer again. However it is possible to divide any two integers with a remainder. So for example 7 divided by 3 is equal to 2 with a remainder of 1, since $7 = 2 \cdot 3 + 1$.

Doing integer division like this is probably something many of us remember from school. It is usually called **Euclidean division**, or **division with a remainder**, and it is an essential technique to understand many concepts in this book. The precise definition is as follows:

Let $a \in \mathbb{Z}$ and $b \in \mathbb{Z}$ be two integers with $b \neq 0$. Then there is always another integer $m \in \mathbb{Z}$ and a counting number $r \in \mathbb{N}$, with $0 \le r < |b|$ such that

$$a = m \cdot b + r \tag{3.7}$$

This decomposition of a given b is called **Euclidean division**, where a is called the **dividend**, b is called the **divisor**, m is called the **quotient** and r is called the **remainder**.

Notation and Symbols 1. Suppose that the numbers a,b,m and r satisfy equation (3.7). Then we often write

$$a \operatorname{div} b := m, \quad a \operatorname{mod} b := r \tag{3.8}$$

to describe the quotient and the remainder of the Euclidean division. We also say that an integer a is divisible by another integer b if a mod b = 0 holds. In this case we also write b|a.

So, in a nutshell Euclidean division is a process of dividing one integer by another in a way that produces a quotient and a non-negative remainder, the latter of which is smaller than the absolute value of the divisor. It can be shown that both the quotient and the remainder always exist and are unique, as long as the dividend is different from 0.

A special situation occurs whenever the remainder is zero, because in this case the dividend is divisible by the divisor. Our notation b|a reflects that.

Example 2. Applying Euclidean division and our previously defined notation 3.8 to the divisor -17 and the dividend 4, we get

```
-17 \text{ div } 4 = -5, \quad -17 \text{ mod } 4 = 3
```

because -17 = -5.4 + 3 is the Euclidean division of -17 and 4 (the remainder is, by definition, a non-negative number). In this case 4 does not divide -17, as the reminder is not zero. The truth value of the expression 4|-17 therefore is FALSE. On the other hand, the truth value of 4|12 is TRUE, since 4 divides 12, as $12 \mod 4 = 0$. We can invoke SageMath to do the computation for us. We get the following:

```
sage: ZZ(-17) // ZZ(4) # Integer quotient
                                                                                22
1042
                                                                                23
1043
    sage: ZZ(-17) % ZZ(4) # remainder
                                                                                24
1044
                                                                                25
1045
    sage: ZZ(4).divides(ZZ(-17)) # self divides other
                                                                                26
1046
    False
                                                                                27
1047
    sage: ZZ(4).divides(ZZ(12))
                                                                                28
1048
    True
                                                                                29
1049
```

Methods to compute Euclidean division for integers are called **integer division algorithms**. Probably the best known algorithm is the so-called **long division**, which most of us might have learned in school. (It should be noted, however that there are faster methods like **Newton–Raphson division**.)

As long division is the standard method used for pen-&-paper division of multi-digit numbers expressed in decimal notation, the reader should become familiar with it as we use it throughout this book when we do simple pen-and-paper computations. However, instead of defining the algorithm formally, we rather give some examples that will hopefully make the process clear.

In a nutshell, the algorithm loops through the digits of the dividend from the left to right, subtracting the largest possible multiple of the divisor (at the digit level) at each stage; the multiples then become the digits of the quotient, and the remainder is the first digit of the dividend.

Example 3 (Integer Long Division). To give an example of integer long division algorithm, lets divide the integer a=143785 by the number b=17. Our goal is therefore to find solutions to equation 3.7, that is, we need to find the quotient $m \in \mathbb{Z}$ and the remainder $r \in \mathbb{N}$ such that $143785 = m \cdot 17 + r$. Using a notation that is mostly used in Commonwealth countries, we

1067 compute as follows:

$$\begin{array}{r}
8457 \\
17)143785 \\
\underline{136} \\
77 \\
\underline{68} \\
98 \\
\underline{85} \\
135 \\
\underline{119} \\
16
\end{array}$$
(3.9)

We therefore get m = 8457 as well as r = 16 and indeed we have $143785 = 8457 \cdot 17 + 16$, which we can double check invoking Sage:

Exercise 5 (Integer Long Division). Find an $m \in \mathbb{Z}$ as well as an $r \in \mathbb{N}$ such that $a = m \cdot b + r$ holds for the following pairs (a,b) = (27,5), (a,b) = (27,-5), (a,b) = (127,0), (a,b) = (-1687,11) and . In which cases are your solutions unique?

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$$(a,b) = (0,7)$$

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Exercise 6 (Long Division Algorithm). Write an algorithm in pseudocode that computes integer pseudocode long division, handling all edge cases properly.

The Extended Euclidean Algorithm One of the most critical parts in this book is modular arithmetics XXX and its application in the computations in so-called **finite fields**, as we explain in section 4.3. In modular arithmetics, it is sometimes possible to define actual division and multiplicative inverses of numbers that is very different from inverses as we know them from other systems like factional numbers.

However, to actually compute those inverses, we have to get familiar with the so-called **extended Euclidean algorithm**. A few more terms are necessary to explain the concept: The **greatest common divisor** (GCD) of two nonzero integers a and b is the greatest non-zero counting number d such that d divides both a and b, that is, d|a as well as d|b. We write gcd(a,b) := d for this number. In addition, two counting numbers are called **relative primes** or **coprimes**, if their greatest common divisor is 1.

The extended Euclidean algorithm is a method to calculate the greatest common divisor of two counting numbers a and $b \in \mathbb{N}$, as well as two additional integers $s,t \in \mathbb{Z}$, such that the following equation holds:

$$gcd(a,b) = s \cdot a + t \cdot b \tag{3.10}$$

The following pseudocode shows in detail how to calculate these numbers with the extended Euclidean algorithm:

The algorithm is simple enough to be done effectively in pen-&-paper examples, where it is common to write it as a table where the rows represent the while-loop and the columns represent the values of the the array r, s and t with index k. The following example provides a simple execution:

check reference

modular arithmetics

actual division

multiplicative inverses

factional numbers

Algorithm 1 Extended Euclidean Algorithm

```
Require: a, b \in \mathbb{N} with a > b
   procedure EXT-EUCLID(a,b)
         r_0 \leftarrow a
         r_1 \leftarrow b
         s_0 \leftarrow 1
         s_1 \leftarrow 0
         k \leftarrow 1
         while r_k \neq 0 do
               q_k \leftarrow r_{k-1} \text{ div } r_k
               r_{k+1} \leftarrow r_{k-1} - q_k \cdot r_k
               s_{k+1} \leftarrow s_{k-1} - q_k \cdot s_k
               k \leftarrow k+1
         end while
         return gcd(a,b) \leftarrow r_{k-1}, s \leftarrow s_{k-1} and t := (r_{k-1} - s_{k-1} \cdot a) div b
   end procedure
Ensure: gcd(a,b) = s \cdot a + t \cdot b
```

Example 4. To illustrate the algorithm, let's apply it to the numbers a = 12 and b = 5. Since $12, 5 \in \mathbb{N}$ as well as $12 \ge 5$ all requirements are meat and we compute as follows:

From this we can see that 12 and 5 are relatively prime (coprime), since their greatest common divisor is gcd(12,5) = 1 and that the equation $1 = (-2) \cdot 12 + 5 \cdot 5$ holds. We can also invoke sage to double check our findings:

Exercise 7 (Extended Euclidean Algorithm). Find integers $s,t \in \mathbb{Z}$ such that $gcd(a,b) = s \cdot a + t \cdot b$ holds for the following pairs (a,b) = (45,10), (a,b) = (13,11), (a,b) = (13,12). What pairs (a,b) are coprime?

Exercise 8 (Towards Prime fields). Let $n \in \mathbb{N}$ be a counting number and p a prime number, such that n < p. What is the greatest common divisor gcd(p,n)?

Exercise 9. Find all numbers $k \in \mathbb{N}$ with $0 \le k \le 100$ such that gcd(100, k) = 5.

Exercise 10. Show that gcd(n,m) = gcd(n+m,m) for all $n,m \in \mathbb{N}$.

3.3 Modular arithmetic

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In mathematics, **modular arithmetic** is a system of arithmetic for integers, where numbers "wrap around" when reaching a certain value, much like calculations on a clock wrap around whenever the value exceeds the number 12. For example, if the clock shows that it is 11 o'clock,

then 20 hours later it will be 7 o'clock, not 31 o'clock. The number 31 has no meaning on a normal clock that shows hours.

The number at which the wrap occurs is called the **modulus**. Modular arithmetics generalizes the clock example to arbitrary moduli and studies equations and phenomena that arize in this new kind of arithmetics. It is of central importance for understanding most modern crypto systems, in large parts because the exponentiation function has an inverse with respect to certain moduli that is hard to compute. In addition, we will see that it provides the foundation of what is called finite fields ().

exponentiation

See XXX

Although modular arithmetic appears very different from ordinary integer arithmetic that we are all familiar with, we encourage the interested reader to work through the example and to discover that, once they get used to the idea that this is a new kind of calculations, it will seem much less daunting.

Congurency In what follows, let $n \in \mathbb{N}$ with $n \ge 2$ be a fixed counting number that we will call the **modulus** of our modular arithmetics system. With such an n given, we can then group integers into classes, by saying that two integers are in the same class, whenever their Euclidean division 3.2 by n will give the same remainder. We then say that two numbers are **congruent** whenever they are in the same class.

Example 5. If we choose n = 12 as in our clock example, then the integers -7, 5, 17 and 29 are all congruent with respect to 12, since all of them have the remainder 5 if we perform Euclidean division on them by 12. In the picture of an analog 12-hour clock, starting at 5 o'clock, when we add 12 hours we are again at 5 o'clock, representing the number 17. On the other hand, when we subtract 12 hours, we are at 5 o'clock again, representing the number -7.

We can formalize this intuition of what congruency should be into a proper definition utilizing Euclidean division (as explained previously in 3.2): Let $a, b \in \mathbb{Z}$ be two integers and $n \in \mathbb{N}$ a natural number. Then a and b are said to be **congruent with respect to the modulus** n, if and only if the following equation holds

$$a \bmod n = b \bmod n \tag{3.11}$$

If, on the other hand, two numbers are not congruent with respect to a given modulus n, we call them **incongruent** w.r.t. n.

A **congruency** is then nothing but an equation "up to congruency", which means that the equation only needs to hold if we take the modulus on both sides. In which case we write

$$a \equiv b \pmod{n} \tag{3.12}$$

Exercise 11. Which of the following pairs of numbers are congruent with respect to the modulus 13: (5,19), (13,0), (-4,9), (0,0).

Exercise 12. Find all integers x, such that the congruency $x \equiv 4 \pmod{6}$ is satisfied.

Modular Arithmetics One particulary useful thing about congruencies is that we can do calculations (arithmetics), much like we can with integer equations. That is, we can add or multiply numbers on both sides. The main difference is probably that the congruency $a \equiv b \pmod{n}$ is only equivalent to the congruency $k \cdot a \equiv k \cdot b \pmod{n}$ for some non zero integer $k \in \mathbb{Z}$, whenever k and the modulus n are coprime. The following list gives a set of useful rules:

Suppose that the congurencies $a_1 \equiv b_1 \pmod{n}$ as well as $a_2 \equiv b_2 \pmod{n}$ are satisfied for integers $a_1, a_2, b_1, b_2 \in \mathbb{Z}$ and that $k \in \mathbb{Z}$ is another integer. Then:

- $a_1 + k \equiv b_1 + k \pmod{n}$ (compatibility with translation)
- $k \cdot a_1 \equiv k \cdot b_1 \pmod{n}$ (compatibility with scaling)
- $a_1 + a_2 \equiv b_1 + b_2 \pmod{n}$ (compatibility with addition)
- $a_1 \cdot a_2 \equiv b_1 \cdot b_2 \pmod{n}$ (compatibility with multiplication)

Other rules, such as compatibility with subtraction and exponentiation, follow from the rules above. For example, compatibility with subtraction follows from compatibility with scaling by k = -1 and compatibility with addition.

Note that the previous rules are implications, not equivalences, which means that you can not necessarily reverse those rules. The following rules makes this precise:

- If $a_1 + k \equiv b_1 + k \pmod{n}$, then $a_1 \equiv b_1 \pmod{n}$
 - If $k \cdot a_1 \equiv k \cdot b_1 \pmod{n}$ and k is coprime with n, then $a_1 \equiv b_1 \pmod{n}$
- If $k \cdot a_1 \equiv k \cdot b_1 \pmod{k \cdot n}$, then $a_1 \equiv b_1 \pmod{n}$

Another property of congruencies, not known in the traditional arithmetics of integers is the Fermat's Little Theorem. In simple words, it states that, in modular arithmetics, every number raised to the power of a prime number modulus is congruent to the number itself. Or, to be more precise, if $p \in \mathbb{P}$ is a prime number and $k \in \mathbb{Z}$ is an integer, then:

$$k^p \equiv k \pmod{p} \,, \tag{3.13}$$

If k is coprime to p, then we can divide both sides of this congruency by k and rewrite the expression into the equivalent form

$$k^{p-1} \equiv 1 \pmod{p} \tag{3.14}$$

This Sage

We can use Sage to compute examples for both k being coprime and not coprime to p:

```
snippet
    sage: ZZ(137).gcd(ZZ(64))
1179
                                                                                    should be
1180
                                                                                    described
    sage: ZZ(64) ** ZZ(137) % ZZ(137) == ZZ(64) % ZZ(137)
1181
                                                                                    <sup>8</sup>n more
    True
                                                                                    detail.
1182
    sage: ZZ(64) ** ZZ(137-1) % ZZ(137) == ZZ(1) % ZZ(137)
                                                                                   40
1183
    True
                                                                                   41
1184
    sage: ZZ(1918).gcd(ZZ(137))
                                                                                   42
1185
    137
                                                                                   43
1186
    sage: ZZ(1918) ** ZZ(137) % ZZ(137) == ZZ(1918) % ZZ(137)
                                                                                   44
1187
    True
                                                                                   45
1188
    sage: ZZ(1918) ** ZZ(137-1) % ZZ(137) == ZZ(1) % ZZ(137)
                                                                                   46
1189
    False
                                                                                   47
1190
```

Now, since for the sake of readers who have never encountered modular arithmetics before, let's compute an example that contains most of the concepts described in this section:

Example 6. Assume that we choose the modulus 6 and that our task is to solve the following congruency for $x \in \mathbb{Z}$

$$7 \cdot (2x+21) + 11 \equiv x - 102 \pmod{6}$$

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As many rules for congruencies are more or less same as for integers, we can proceed in a similar way as we would if we had an equation to solve. Since both sides of a congruency contain ordinary integers, we can rewrite the left side as follows: $7 \cdot (2x + 21) + 11 = 14x + 147 = 14x + 158$. We can therefore rewrite the congruency into the equivalent form

$$14x + 158 \equiv x - 102 \pmod{6}$$

In the next step we want to shift all instances of x to left and every other term to the right. So we apply the "compatibility with translation" rules two times. In a first step we choose k = -x and in a second step we choose k = -158. Since "compatibility with translation" transforms a congruency into an equivalent form, the solution set will not change and we get

$$14x + 158 \equiv x - 102 \pmod{6} \Leftrightarrow 14x - x + 158 - 158 \equiv x - x - 102 - 158 \pmod{6} \Leftrightarrow 13x \equiv -260 \pmod{6}$$

If our congruency would just be a normal integer equation, we would divide both sides by 13 to get x = -20 as our solution. However, in case of a congruency, we need to make sure that the modulus and the number we want to divide by are coprime first – only then will we get an equivalent expression. So we need to find the greatest common divisor gcd(13,6). Since 13 is prime and 6 is not a multiple of 13, we know that gcd(13,6) = 1, so these numbers are indeed coprime. We therefore compute

$$13x \equiv -260 \pmod{6} \Leftrightarrow x \equiv -20 \pmod{6}$$

Our task is now to find all integers x, such that x is congruent to -20 with respect to the modulus 6. So we have to find all x such

$$x \mod 6 = -20 \mod 6$$

Since $-4 \cdot 6 + 4 = -20$ we know $-20 \mod 6 = 4$ and hence we know that x = 4 is a solution to this congruency. However, 22 is another solution since 22 mod 6 = 4 as well, and so is -20. In fact, there are infinitely many solutions given by the set

$$\{\ldots, -8, -2, 4, 10, 16, \ldots\} = \{4 + k \cdot 6 \mid k \in \mathbb{Z}\}$$

Putting all this together, we have shown that the every x from the set $\{x = 4 + k \cdot 6 \mid k \in \mathbb{Z}\}$ is a solution to the congruency $7 \cdot (2x + 21) + 11 \equiv x - 102 \pmod{6}$. We double ckeck for, say, x = 4 as well as $x = 14 + 12 \cdot 6 = 86$ using sage:

```
sage: (ZZ(7)*(ZZ(2)*ZZ(4) + ZZ(21)) + ZZ(11)) % ZZ(6) == (ZZ(2)) % (ZZ(2
1196
                                                                              (4) - ZZ(102)
                                                                                                                                                                                                                                                               % ZZ(6)
1197
                                           True
                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                    49
1198
                                           sage: (ZZ(7)*(ZZ(2)*ZZ(76) + ZZ(21)) + ZZ(11)) % ZZ(6) == (
                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                    50
1199
                                                                           ZZ(76) - ZZ(102)
                                                                                                                                                                                                                                                                                                 % ZZ(6)
1200
                                           True
                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                                    51
1201
```

Readers who had not been familiar with modular arithmetics until now and who might be discouraged by how complicated modular arithmetics seems at this point, should keep two things in mind. First, computing congruencies in modular arithmetics is not really more complicated than computations in more familiar number systems (e.g. fractional numbers), it is just

a matter of getting used to it. Second, the theory of prime fields (and more general residue class rings) takes a different view on modular arithmetics with the attempt to simplify matters. In other words, once we understand prime field arithmetics, things become conceptually cleaner and more easy to compute.

ields

class rings

Exercise 13. Choose the modulus 13 and find all solutions $x \in \mathbb{Z}$ to the following congruency

 $5x + 4 \equiv 28 + 2x \pmod{13}$

Exercise 14. Choose the modulus 23 and find all solutions $x \in \mathbb{Z}$ to the following congruency $69x \equiv 5 \pmod{23}$

Exercise 15. Choose the modulus 23 and find all solutions $x \in \mathbb{Z}$ to the following congruency $69x \equiv 46 \pmod{23}$

The Chinese Remainder Theorem We have seen how to solve congruencies in modular arithmetic. However, one question that remains is how to solve systems of congruencies with different moduli? The answer is given by the Chinese reimainder theorem, which states that for any $k \in \mathbb{N}$ and coprime natural numbers $n_1, \ldots n_k \in \mathbb{N}$ as well as integers $a_1, \ldots a_k \in \mathbb{Z}$, the so-called simultaneous congruency

$$x \equiv a_1 \pmod{n_1}$$

$$x \equiv a_2 \pmod{n_2}$$

$$\dots$$

$$x \equiv a_k \pmod{n_k}$$
(3.15)

has a solution, and all possible solutions of this congruence system are congruent modulo the product $N = n_1 \cdot ... \cdot n_k$. In fact, the following algorithm computes the solution set:

check algorithm floating

Algorithm 2 Chinese Remainder Theorem

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```
Require: n_0, \ldots, n_{k-1} \in \mathbb{N} coprime 

procedure CONGRUENCY-SYSTEMS-SOLVER(k, a_{0, \ldots, k-1}, n_{0, \ldots, k-1}) N \leftarrow n_0 \cdot \ldots \cdot n_{k-1} while j < k do  N_j \leftarrow N/n_j \\  (\_, s_j, t_j) \leftarrow EXT - EUCLID(N_j, n_j)  \Rightarrow 1 = s_j \cdot N_j + t_j \cdot n_j end while  x' \leftarrow \sum_{j=0}^{k-1} a_j \cdot s_j \cdot N_j \\  x \leftarrow x' \bmod N \\  \mathbf{return} \ \{x + m \cdot N \mid m \in \mathbb{Z}\}  end procedure
```

Example 7. To illustrate how to solve simultaneous congruences using the Chinese remainder theorem, let's look at the following system of congruencies:

Ensure: $\{x + m \cdot N \mid m \in \mathbb{Z}\}$ is the complete solution set to 3.15.

$$x \equiv 4 \pmod{7}$$

$$x \equiv 1 \pmod{3}$$

$$x \equiv 3 \pmod{5}$$

$$x \equiv 0 \pmod{11}$$

¹This is the classical Chinese remainder theorem as it was already known in ancient China. Under certain circumstances, the theorem can be extended to non-coprime moduli n_1, \ldots, n_k but this is beyond the scope of this book. Interested readers should consult XXX add references

Clearly all moduli are coprime and we have $N = 7 \cdot 3 \cdot 5 \cdot 11 = 1155$, as well as $N_1 = 165$, $N_2 = 385$, $N_3 = 231$ and $N_4 = 105$. From this we calculate with the extended Euclidean algorithm

$$1 = 2 \cdot 165 + -47 \cdot 7$$

$$1 = 1 \cdot 385 + -128 \cdot 3$$

$$1 = 1 \cdot 231 + -46 \cdot 5$$

$$1 = 2 \cdot 105 + -19 \cdot 11$$

so we have $x = 4 \cdot 2 \cdot 165 + 1 \cdot 1 \cdot 385 + 3 \cdot 1 \cdot 231 + 0 \cdot 2 \cdot 105 = 2398$ as one solution. Because 2398 mod 1155 = 88 the set of all solutions is $\{\dots, -2222, -1067, 88, 1243, 2398, \dots\}$. In particular, there are infinitely many different solutions. We can invoke Sage's computation of the Chinese Remainder Theorem (CRT) to double check our findings:

As we have seen in various examples before, computing congruencies can be cumbersome and solution sets are large in general. It is therefore advantaegous to find some kind of simplification for modular arithmetic.

Fortunately, this is possible and relatively straightforward once we consider all integers that have the same remainder with respect to a given modulus n in Euclidean division to be equivalent. Then we can go a step further, and identify each set of numbers with equal remainder with that remainder and call it a **remainder class** or **residue class** in modulo n arithmetics.

It then follows from the properties of Euclidean division that there are exactly n different remainder classes for every modulus n and that integer addition and multiplication can be projected to a new kind of addition and multiplication on those classes.

Roughly speaking, the new rules for addition and multiplication are then computed by taking any element of the first equivalence class and some element of the second, then add or multiply them in the usual way and see which equivalence class the result is contained in. The following example makes this abstract description more concrete:

Example 8 (Arithmetics modulo 6). Choosing the modulus n = 6, we have six equivalence classes of integers which are congruent modulo 6 (they have the same remainder when divided by 6) and when we identify each of those remainder classes with the remainder, we get the following identification:

$$0 := \{..., -6, 0, 6, 12, ...\}$$

$$1 := \{..., -5, 1, 7, 13, ...\}$$

$$2 := \{..., -4, 2, 8, 14, ...\}$$

$$3 := \{..., -3, 3, 9, 15, ...\}$$

$$4 := \{..., -2, 4, 10, 16, ...\}$$

$$5 := \{..., -1, 5, 11, 17, ...\}$$

Now to compute the addition of those equivalence classes, say 2+5, one chooses arbitrary elements from both sets, say 14 and -1, adds those numbers in the usual way and then looks at the equivalence class of the result.

So we get 14 + (-1) = 13, and 13 is in the equivalence class (of) 1. Hence we find that 2 + 5 = 1 in modular 6 arithmetics, which is a more readable way to write the congruency $2 + 5 \equiv 1 \pmod{6}$.

Applying the same reasoning to all equivalence classes, addition and multiplication can be transferred to equivalence classes. The results for modulus 6 arithmetics are summarized in the following addition and multiplication tables:

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+	0	1	2	3	4	5		•	0	1	2	3	4	5
0	0	1	2	3	4	5	•	0	0	0	0	0	0	0
1	1	2	3	4	5	0		1	0	1	2	3	4	5
2	2	3	4	5	0	1		2	0	2	4	0	2	4
3	3	4	5	0	1	2		3	0	3	0	3	0	3
4	4	5	0	1	2	3		4	0	4	2	0	4	2
5	5	0	1	2	3	4		5	0	5	2	3	2	1

This way, we have defined a new arithmetic system that contains just 6 numbers and comes with its own definition of addition and multiplication. It is called **modular 6 arithmetics** and written as \mathbb{Z}_6 .

To see why such an identification of a congruency class with its remainder is useful and actually simplifies congruency computations a lot, lets go back to the congruency from example 6 again:

$$7 \cdot (2x+21) + 11 \equiv x - 102 \pmod{6}$$
 (3.16)

As shown in example 6, the arithmetics of congruencies can deviate from ordinary arithmetics: For example, division needs to check whether the modulus and the dividend are coprimes, and solutions are not unique in general.

We can rewrite this congruency as an **equation** over our new arithmetic type \mathbb{Z}_6 by **projecting onto the remainder classes**. In particular, since $7 \mod 6 = 1$, $21 \mod 6 = 3$, $11 \mod 6 = 5$ and $102 \mod 6 = 0$ we have

$$7 \cdot (2x+21) + 11 \equiv x - 102 \pmod{6}$$
 over \mathbb{Z}
 $\Leftrightarrow 1 \cdot (2x+3) + 5 = x$ over \mathbb{Z}_6

We can use the multiplication and addition table above to solves the equation on the right like we would solve normal integer equations:

$$1 \cdot (2x+3) + 5 = x$$

$$2x+3+5 = x$$

$$2x+2 = x$$

$$2x+2+4-x = x+4-x$$

$$x = 4$$
addition-table: $2+4=0$

As we can see, despite the somewhat unfamiliar rules of addition and multiplication, solving congruencies this way is very similar to solving normal equations. And, indeed, the solution set is identical to the solution set of the original congruency, since 4 is identified with the set $\{4+6\cdot k\mid k\in\mathbb{Z}\}$.

We can invoke Sage to do computations in our modular 6 arithmetics type. This is particularly useful to double-check our computations:

```
1270 sage: Z6 = Integers(6) 54

1271 sage: Z6(2) + Z6(5) 55

1272 1 56

1273 sage: Z6(7)*(Z6(2)*Z6(4)+Z6(21))+Z6(11) == Z6(4) - Z6(102) 57

1274 True 58
```

Jargon 1 (k-bit modulus). In cryptographic papers, we can sometimes read phrases like"[...] using a 4096-bit modulus". This means that the underlying modulus n of the modular arithmetic used in the system has a binary representation with a length of 4096 bits. In contrast, the number 6 has the binary representation 110 and hence our example 8 describes a 3-bit modulus arithmetics system.

Exercise 16. Let a, b, k be integers, such that $a \equiv b \pmod{n}$ holds. Show $a^k \equiv b^k \pmod{n}$.

Exercise 17. Let a, n be integers, such that a and n are not coprime. For which $b \in \mathbb{Z}$ does the congruency $a \cdot x \equiv b \pmod{n}$ have a solution x and how does the solution set look in that case?

Modular Inverses As we know, integers can be added, subtracted and multiplied so that the result is also an integer, but this is not true for the division of integers in general: for example, 3/2 is not an integer anymore. To see why this is, from a more theoretical perspective, let us consider the definition of a multiplicative inverse first. When we have a set that has some kind of multiplication defined on it and we have a distinguished element of that set that behaves neutrally with respect to that multiplication (doesn't change anything when multiplied with any other element), then we can define **multiplicative inverses** in the following way:

Let *S* be our set that has some notion $a \cdot b$ of multiplication and a **neutral element** $1 \in S$, such that $1 \cdot a = a$ for all elements $a \in S$. Then a **multiplicative inverse** a^{-1} of an element $a \in S$ is defined as follows:

$$a \cdot a^{-1} = 1 \tag{3.17}$$

Informally speaking, the definition of a multiplicative inverse is means that it "cancels" the original element to give 1 when they are multiplied.

Numbers that have multiplicative inverses are of particular interest, because they immediately lead to the definition of division by those numbers. In fact, if a is number such that the multiplicative inverse a^{-1} exists, then we define **division** by a simply as multiplication by the inverse:

$$\frac{b}{a} := b \cdot a^{-1} \tag{3.18}$$

Example 9. Consider the set of rational numbers, also known as fractions, \mathbb{Q} . For this set, the neural element of multiplication is 1, since $1 \cdot a = a$ for all rational numbers. For example, $1 \cdot 4 = 4$, $1 \cdot \frac{1}{4} = \frac{1}{4}$, or $1 \cdot 0 = 0$ and so on.

Every rational number $a \neq 0$ has a multiplicative inverse, given by $\frac{1}{a}$. For example, the multiplicative inverse of 3 is $\frac{1}{3}$, since $3 \cdot \frac{1}{3} = 1$, the multiplicative inverse of $\frac{5}{7}$ is $\frac{7}{5}$, since $\frac{5}{7} \cdot \frac{7}{5} = 1$, and so on

Example 10. Looking at the set \mathbb{Z} of integers, we see that with respect to multiplication the neutral element is the number 1 and we notice that no integer $a \neq 1$ has a multiplicative inverse, since the equation $a \cdot x = 1$ has no integer solutions for $a \neq 1$.

The definition of multiplicative inverse works verabtim for addition as well. In the case of integers, the neutral element with respect to addition is 0, since a + 0 = 0 for all integers $a \in \mathbb{Z}$. The additive inverse always exist and is given by the negative number -a, since a + (-1) = 0.

Example 11. Looking at the set \mathbb{Z}_6 of residual classes modulo 6 from example 8, we can use the multiplication table to find multiplicative inverses. To do so, we look at the row of the element and then find the entry equal to 1. If such an entry exists, the element of that column is the multiplicative inverse. If, on the other hand, the row has no entry equal to 1, we know that the element has no multiplicative inverse.

(-1) should be (-a)?

For example in \mathbb{Z}_6 the multiplicative inverse of 5 is 5 itself, since $5 \cdot 5 = 1$. We can also see that 5 and 1 are the only elements that have multiplicative inverses in \mathbb{Z}_6 .

Now, since 5 has a multiplicative inverse modulo 6, it makes sense to "divide" by 5 in \mathbb{Z}_6 . For example

$$\frac{4}{5} = 4 \cdot 5^{-1} = 4 \cdot 5 = 2$$

From the last example, we can make the interesting observation that while 5 has no multiplicative inverse as an integer, it has a multiplicative inverse in modular 6 arithmetics.

The remaining question is to understand which elements have multiplicative inverses in modular arithmetics. The answer is that, in modular n arithmetics, a residue class r has a multiplicative inverse, if and only if n and r are coprime. Since ggt(n,r)=1 in that case, we know from the extended Euclidean algorithm that there are numbers s and t, such that

$$1 = s \cdot n + t \cdot r \tag{3.19}$$

If we take the modulus n on both sides, the term $s \cdot n$ vanishes, which tells us that $t \mod n$ is the multiplicative inverse of r in modular n arithmetics.

Example 12 (Multiplicative inverses in \mathbb{Z}_6). In the previous example, we looked up multiplicative inverses in \mathbb{Z}_6 from the lookup-table in Example 8. In real world examples, it is usually impossible to write down those lookup tables, as the modulus is way too large, and the sets occasionally contain more elements than there are atoms in the observable universe.

Now, trying to determine that $2 \in \mathbb{Z}_6$ has no multiplicative inverse in \mathbb{Z}_6 without using the lookup table, we immediately observe that 2 and 6 are not coprime, since their greatest common divisor is 2. It follows that equation 3.19 has no solutions s and t, which means that 2 has no multiplicative inverse in \mathbb{Z}_6 .

The same reasoning works for 3 and 4, as neither of these are coprime with 6. The case of 5 is different, since ggt(6,5) = 1. To compute the multiplicative inverse of 5, we use the extended Euclidean algorithm and compute the following:

We get s = 1 as well as t = -1 and have $1 = 1 \cdot 6 - 1 \cdot 5$. From this, it follows that $-1 \mod 6 = 5$ is the multiplicative inverse of 5 in modular 6 arithmetics. We can double check using Sage:

At this point, the attentive reader might notice that the situation where the modulus is a prime number is of particular interest, because we know from exercise XXX that in these cases all remainder classes must have modular inverses, since ggt(r,n) = 1 for prime n and r < n. In fact, Fermat's little theorem provides a way to compute multiplicative inverses in this situation, since in case of a prime modulus p and r < p, we get the following:

$$r^p \equiv r \pmod{p} \Leftrightarrow$$
 $r^{p-1} \equiv 1 \pmod{p} \Leftrightarrow$
 $r \cdot r^{p-2} \equiv 1 \pmod{p}$

This tells us that the multiplicative inverse of a residue class r in modular p arithmetic is precisely r^{p-2} .

Example 13 (Modular 5 arithmetics). To see the unique properties of modular arithmetics whenever the modulus is a prime number, we will replicate our findings from example 8, but this time for the prime modulus 5. For n = 5 we have five equivalence classes of integers which are congruent modulo 5. We write this as follows:

$$0 := \{\dots, -5, 0, 5, 10, \dots\}$$

$$1 := \{\dots, -4, 1, 6, 11, \dots\}$$

$$2 := \{\dots, -3, 2, 7, 12, \dots\}$$

$$3 := \{\dots, -2, 3, 8, 13, \dots\}$$

$$4 := \{\dots, -1, 4, 9, 14, \dots\}$$

Addition and multiplication can be transferred to the equivalence classes, in a way exactly parallel to Example 8. This results in the following addition and multiplication tables:

+	0	1	2	3	4	•	0	1	2	3	4
0	0	1	2	3	4	0	0	0	0	0	0
1	1	2	3	4	0	1	0	1	2	3	4
2	2	3	4	0	1	2	0	2	4	1	3
3	3	4	0	1	2	3	0	3	1	4	2
4	4	0	1	2	3	4	0	4	3	2	1

Calling the set of remainder classes in modular 5 arithmetics with this addition and multiplication \mathbb{F}_5 (for reasons we explain in more detail in XXX), we see some subtle but important differences to the situation in \mathbb{Z}_6 . In particular, we see that in the multiplication table, every remainder $r \neq 0$ has the entry 1 in its row and therefore has a multiplicative inverse. In addition, there are no non-zero elements such that their product is zero.

add reference

To use Fermat's little theorem in \mathbb{F}_5 for computing multiplicative inverses (instead of using the multiplication table), lets consider $3 \in \mathbb{F}_3$. We know that the multiplicative inverse is given by the remainder class that contains $3^{5-2} = 3^3 = 3 \cdot 3 \cdot 3 = 4 \cdot 3 = 2$. And indeed $3^{-1} = 2$, since $3 \cdot 2 = 1$ in \mathbb{F}_5 .

We can invoke Sage to do computations in our modular 5 arithmetics type to double-check our computations:

```
sage: Z5 = Integers(5)
                                                                                        61
1359
     sage: Z5(3) ** (5-2)
                                                                                        62
1360
                                                                                        63
1361
     sage: Z5(3) ** (-1)
                                                                                        64
1362
                                                                                        65
1363
     sage: Z5(3)**(5-2) == Z5(3)**(-1)
                                                                                        66
1364
     True
                                                                                        67
1365
```

Example 14. To understand one of the principal differences between prime number modular arithmetics and non-prime number modular arithmetics, consider the linear equation $a \cdot x + b = 0$ defined over both types \mathbb{F}_5 and \mathbb{Z}_6 . Since in \mathbb{F}_5 every non zero element has a multiplicative inverse, we can always solve these types of equations in \mathbb{F}_5 , which is not true in \mathbb{Z}_6 . To see that,

consider the equation 3x + 3 = 0. In \mathbb{F}_5 we have the following:

$$3x+3=0$$
 # add 2 and on both sides
 $3x+3+2=2$ # addition-table: $2+3=0$
 $3x=2$ # divide by 3
 $2\cdot(3x)=2\cdot 2$ # multiplication-table: $2+2=4$
 $x=4$

So in the case of our prime number modular arithmetics, we get the unique solution x = 4. Now consider \mathbb{Z}_6 :

$$3x + 3 = 0$$
 # add 3 and on both sides
 $3x + 3 + 3 = 3$ # addition-table: $3 + 3 = 0$
 $3x = 3$ # no multiplicative inverse of 3 exists

So, in this case, we cannot solve the equation for x by dividing by 3. And, indeed, when we look at the multiplication table of \mathbb{Z}_6 (Example 8), we find that there are three solutions $x \in \{1,3,5\}$, such that 3x + 3 = 0 holds true for all of them.

Exercise 18. Consider the modulus n = 24. Which of the integers 7, 1, 0, 805, -4255 have multiplicative inverses in modular 24 arithmetics? Compute the inverses, in case they exist.

Exercise 19. Find the set of all solutions to the congruency $17(2x+5)-4 \equiv 2x+4 \pmod{5}$. Then project the congruency into \mathbb{F}_5 and solve the resulting equation in \mathbb{F}_5 . Compare the results.

Exercise 20. Find the set of all solutions to the congruency $17(2x+5)-4 \equiv 2x+4 \pmod{6}$

Then project the congruency into \mathbb{Z}_6 and try to solve the resulting equation in \mathbb{Z}_6 .

3.4 Polynomial Arithmetics

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A polynomial is an expression consisting of variables (also called indeterminates) and coefficients that involves only the operations of addition, subtraction, multiplication, and nonnegative integer exponentiation of variables. All coefficients of a polynomial must have the same type, e.g. being integers or fractions etc. To be more precise a *univariate polynomial* is an expression

$$P(x) := \sum_{j=0}^{m} a_j t^j = a_m x^m + a_{m-1} x^{m-1} + \dots + a_1 x + a_0,$$
 (3.20)

where x is called the **indeterminate**, each a_j is called a **coefficient**. If R is the type of the coefficients, then the set of all **univariate polynomials with coefficients in** R is written as R[x]. We often simply use **polynomial** instead of univariate polynomial, write $P(x) \in R[x]$ for a polynomial and denote the constant term as P(0).

A polynomial is called the **zero polynomial** if all coefficients are zero and a polynomial is called the **one polynomial** if the constant term is 1 and all other coefficients are zero.

If an univariate polynomial $P(x) = \sum_{j=0}^{m} a_j x^j$ is given that is not the zero polynomial, we call deg(P) := m the *degree* of P and define the degree of the zero polynomial to be $-\infty$, where $-\infty$ (negative infinity) is a symbol with the property that $-\infty + m = -\infty$ for all counting numbers $m \in \mathbb{N}$. In addition, we write

$$Lc(P) := a_m \tag{3.21}$$

and call it the **leading coefficient** of the polynomial P. We can restrict the set R[x] of **all** polynomials with coefficients in R, to the set of all such polynomials that have a degree that does not exceed a certain value. If m is the maximum degree allowed, we write $R_{\leq m}[x]$ for the set of all polynomials with a degree less than or equal to m.

Example 15 (Integer Polynomials). The coefficients of a polynomial must all have the same type. The set of polynomials with integer coefficients is written as $\mathbb{Z}[x]$. Examples of such polynomials are:

$$P_1(x) = 2x^2 - 4x + 17$$
 # with $deg(P_1) = 2$ and $Lc(P_1) = 2$
 $P_2(x) = x^{23}$ # with $deg(P_2) = 23$ and $Lc(P_2) = 1$
 $P_3(x) = x$ # with $deg(P_3) = 1$ and $Lc(P_3) = 1$
 $P_4(x) = 174$ # with $deg(P_4) = 0$ and $Lc(P_4) = 174$
 $P_5(x) = 1$ # with $deg(P_5) = 0$ and $Lc(P_5) = 1$
 $P_6(x) = 0$ # with $deg(P_5) = -\infty$ and $Lc(P_6) = 0$
 $P_7(x) = (x - 2)(x + 3)(x - 5)$

In particular, every integer can be seen as an integer polynomial of degree zero. P_7 is a polynomial, because we can expand its definition into $P_7(x) = x^3 - 4x^2 - 11x + 30$, which is polynomial of degree 3 and leading coefficient 1. The following expressions are not integer polynomials:

$$Q_1(x) = 2x^2 + 4 + 3x^{-2}$$

$$Q_2(x) = 0.5x^4 - 2x$$

$$Q_3(x) = 1/x$$

We can invoke Sage to do computations with polynomials. To do so, we have to specify the symbol for the inderteminate and the type for the coefficients. Note, however that Sage defines the degree of the zero polynomial to be -1.

```
sage: Zx = ZZ['x'] # integer polynomials with indeterminate x
                                                                                68
1398
    sage: Zt.<t> = ZZ[] # integer polynomials with indeterminate t
                                                                                69
1399
    sage: Zx
                                                                                70
1400
    Univariate Polynomial Ring in x over Integer Ring
                                                                                71
1401
    sage: Zt
                                                                                72
1402
    Univariate Polynomial Ring in t over Integer Ring
                                                                                73
1403
    sage: p1 = Zx([17,-4,2])
                                                                                74
1404
    sage: p1
                                                                                75
1405
    2*x^2 - 4*x + 17
                                                                                76
1406
    sage: p1.degree()
                                                                                77
1407
                                                                                78
1408
    sage: p1.leading_coefficient()
                                                                                79
1409
                                                                                80
1410
    sage: p2 = Zt(t^23)
                                                                                81
1411
    sage: p2
                                                                                82
1412
    t^23
                                                                                83
1413
    sage: p6 = Zx([0])
                                                                                84
1414
    sage: p6.degree()
                                                                                85
1415
    -1
                                                                                86
1416
```

Example 16 (Polynomials over \mathbb{Z}_6). Recall our definition of the residue classes \mathbb{Z}_6 and their arithmetics as defined in Example 8. The set of all polynomials with indeterminate x and coefficients in \mathbb{Z}_6 is symbolized as $\mathbb{Z}_6[x]$. Example of polynomials from \mathbb{Z}_6 are:

$$P_1(x) = 2x^2 - 4x + 5$$
 # with $deg(P_1) = 2$ and $Lc(P_1) = 2$
 $P_2(x) = x^{23}$ # with $deg(P_2) = 23$ and $Lc(P_2) = 1$
 $P_3(x) = x$ # with $deg(P_3) = 1$ and $Lc(P_3) = 1$
 $P_4(x) = 3$ # with $deg(P_4) = 0$ and $Lc(P_4) = 3$
 $P_5(x) = 1$ # with $deg(P_5) = 0$ and $Lc(P_5) = 1$
 $P_6(x) = 0$ # with $deg(P_5) = -\infty$ and $Lc(P_6) = 0$
 $P_7(x) = (x - 2)(x + 3)(x - 5)$

Just like in the previous example, P_7 is a polynomial. However, since we are working with coefficients from \mathbb{Z}_6 now the expansion of P_7 is computed differently, as we have to invoke addition and multiplication in \mathbb{Z}_6 as defined in XXX. We get the following:

$$(x-2)(x+3)(x-5) = (x+4)(x+3)(x+1)$$

$$= (x^2+4x+3x+3\cdot4)(x+1)$$

$$= (x^2+1x+0)(x+1)$$

$$= (x^3+x^2+x^2+x)$$
additive inverses in \mathbb{Z}_6
bracket expansion
computation in \mathbb{Z}_6
bracket expansion
$$= (x^3+2x^2+x)$$

Again, we can use Sage to do computations with polynomials that have their coefficients in \mathbb{Z}_6 . To do so, we have to specify the symbol for the inderteminate and the type for the coefficients:

```
sage: Z6 = Integers(6)
                                                                               87
1420
    sage: Z6x = Z6['x']
                                                                               88
1421
    sage: Z6x
                                                                               89
1422
    Univariate Polynomial Ring in x over Ring of integers modulo 6
                                                                               90
1423
    sage: p1 = Z6x([5,-4,2])
                                                                               91
1424
    sage: p1
                                                                               92
1425
    2*x^2 + 2*x + 5
                                                                               93
1426
    sage: p1 = Z6x([17,-4,2])
                                                                               94
1427
    sage: p1
                                                                               95
1428
    2*x^2 + 2*x + 5
                                                                               96
    sage: Z6x(x-2)*Z6x(x+3)*Z6x(x-5) == Z6x(x^3 + 2*x^2 + x)
                                                                               97
1430
                                                                               98
1431
```

Given some element from the same type as the coefficients of a polynomial, the polynomial can be evaluated at that element, which means that we insert the given element for every ocurrence of the indeterminate *x* in the polynomial expression.

To be more precise, let $P \in R[x]$, with $P(x) = \sum_{j=0}^{m} a_j x^j$ be a polynomial with a coefficient of type R and let $b \in R$ be an element of that type. Then the **evaluation** of P at b is given as follows:

$$P(a) = \sum_{j=0}^{m} a_j b^j$$
 (3.22)

Example 17. Consider the integer polynomials from example XXX again. To evaluate them at given points, we have to insert the point for all occurences of x in the polynomial expression. Inserting arbitrary values from \mathbb{Z} , we get:

$$P_{1}(2) = 2 \cdot 2^{2} - 4 \cdot 2 + 17 = 17$$

$$P_{2}(3) = 3^{23} = 94143178827$$

$$P_{3}(-4) = -4 = -4$$

$$P_{4}(15) = 174$$

$$P_{5}(0) = 1$$

$$P_{6}(1274) = 0$$

$$P_{7}(-6) = (-6 - 2)(-6 + 3)(-6 + 5) = -264$$

Note, however, that it is not possible to evaluate any of those polynomial on values of different type. For example, it is not strictly correct to write $P_1(0.5)$, since 0.5 is not an integer. We can verify our computations using Sage:

1441 sage:
$$Zx = ZZ['x']$$
 99
1442 sage: $p1 = Zx([17,-4,2])$ 100
1443 sage: $p7 = Zx(x-2)*Zx(x+3)*Zx(x-5)$ 101
1444 sage: $p1(ZZ(2))$ 102
1445 17 103
1446 sage: $p7(ZZ(-6)) == ZZ(-264)$ 104
147 True

Example 18. Consider the polynomials with coefficients in \mathbb{Z}_6 from example XXX again. To evaluate them at given values from \mathbb{Z}_6 , we have to insert the point for all occurences of x in the polynomial expression. We get the following:

$$P_1(2) = 2 \cdot 2^2 - 4 \cdot 2 + 5 = 2 - 2 + 5 = 5$$

$$P_2(3) = 3^{23} = 3$$

$$P_3(-4) = P_3(2) = 2$$

$$P_5(0) = 1$$

$$P_6(4) = 0$$

Exercise 21. Compare both expansions of P_7 from $\mathbb{Z}[x]$ and from $\mathbb{Z}_6[x]$ in example XXX and example XXX, and consider the definition of \mathbb{Z}_6 as given in example XXX. Can you see how the definition of P_7 over \mathbb{Z} projects to the definition over \mathbb{Z}_6 if you consider the residue classes of \mathbb{Z}_6 ?

Polynomial Arithmetics Polynomials behave like integers in many ways. In particular, they can be added, subtracted and multiplied. In addition, they have their own notion of Euclidean division. Informally speaking, we can add two polynomials by simply adding the coefficients of the same index, and we can multiply them by applying the distributive property, that is, by multiplying every term of the left factor with every term of the right factor and adding the results together.

To be more precise let $\sum_{n=0}^{m_1} a_n x^n$ and $\sum_{n=0}^{m_2} b_n x$ be two polynomials from R[x]. Then the **sum** and the **product** of these polynomials is defined as follows:

$$\sum_{n=0}^{m_1} a_n x^n + \sum_{n=0}^{m_2} b_n x^n = \sum_{n=0}^{\max(\{m_1, m_2\})} (a_n + b_n) x^n$$
(3.23)

$$\left(\sum_{n=0}^{m_1} a_n x^n\right) \cdot \left(\sum_{n=0}^{m_2} b_n x^n\right) = \sum_{n=0}^{m_1 + m_2} \sum_{i=0}^n a_i b_{n-i} x^n$$
(3.24)

A rule for polynomial subtraction can be deduced from these two rules by first multiplying the subtrahend with (the polynomial) -1 and then add the result to the minuend.

subtrahend

Regarding the definition of the degree of a polynomial, we see that the degree of the sum is always the maximum of the degrees of both summands, and the degree of the product is always the degree of the factors, since we defined $-\infty \cdot m = \infty$ for every integer $m \in \mathbb{Z}$. Using Sage's definition of degree, this would not hold, as the zero polynomials degree is -1 is Sage, which would violate this rule.

Example 19. To given an example of how polynomial arithmetics works, consider the following two integer polynomials $P, Q \in \mathbb{Z}[x]$ with $P(x) = 5x^2 - 4x + 2$ and $Q(x) = x^3 - 2x^2 + 5$. The sum of these two polynomials is computed by adding the coefficients of each term with equal exponent in x. This gives the following:

$$(P+Q)(x) = (0+1)x^3 + (5-2)x^2 + (-4+0)x + (2+5)$$

= $x^3 + 3x^2 - 4x + 7$

The product of these two polynomials is computed by multiplication of each term in the first factor with each term in the second factor. We get the following:

$$(P \cdot Q)(x) = (5x^2 - 4x + 2) \cdot (x^3 - 2x^2 + 5)$$

= $(5x^5 - 10x^4 + 25x^2) + (-4x^4 + 8x^3 - 20x) + (2x^3 - 4x^2 + 10)$
= $5x^5 - 14x^4 + 10x^3 + 21x^2 - 20x + 10$

1475 sage:
$$Zx = ZZ['x']$$
 111
1476 sage: $P = Zx([2,-4,5])$ 112
1477 sage: $Q = Zx([5,0,-2,1])$ 113
1478 sage: $P+Q == Zx(x^3 +3*x^2 -4*x +7)$ 114
1479 True 115
1480 sage: $P*Q == Zx(5*x^5 -14*x^4 +10*x^3+21*x^2-20*x +10)$ 116
1481 True 117

Example 20. Let us consider the polynomials of the previous example but interpreted in modular 6 arithmetics. So we consider $P, Q \in \mathbb{Z}_6[x]$ again with $P(x) = 5x^2 - 4x + 2$ and $Q(x) = x^3 - 2x^2 + 5$. This time we get the following:

$$(P+Q)(x) = (0+1)x^3 + (5-2)x^2 + (-4+0)x + (2+5)$$
$$= (0+1)x^3 + (5+4)x^2 + (2+0)x + (2+5)$$
$$= x^3 + 3x^2 + 2x + 1$$

$$(P \cdot Q)(x) = (5x^2 - 4x + 2) \cdot (x^3 - 2x^2 + 5)$$

$$= (5x^2 + 2x + 2) \cdot (x^3 + 4x^2 + 5)$$

$$= (5x^5 + 2x^4 + 1x^2) + (2x^4 + 2x^3 + 4x) + (2x^3 + 2x^2 + 4)$$

$$= 5x^5 + 4x^4 + 4x^3 + 3x^2 + 4x + 4$$

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Exercise 22. Compare the sum P+Q and the product $P\cdot Q$ from the previous two examples XXX and XXX and consider the definition of \mathbb{Z}_6 as given in example XXX. How can we derive the computations in $\mathbb{Z}_6[x]$ from the computations in Z[x]?

Euklidean Division The ring of polynomials shares a lot of properties with integers. In particular, the concept of Euclidean division and the algorithm of long division is also defined for polynomials. Recalling the Euclidean division of integers XXX, we know that, given two integers a and $b \neq 0$, there is always another integer a and a counting number a with a such that $a = m \cdot b + r$ holds.

We can generalize this to polynomials whenever the leading coefficient of the dividend polynomial has a notion of multiplicative inverse. In fact, given two polynomials A and $B \neq 0$ from R[x] such that $Lc(B)^{-1}$ exists in R, there exist two polynomials M (the quotient) and R (the remainder), such that the following equation holds:

$$A = M \cdot B + R \tag{3.25}$$

and deg(R) < deg(B). Similarly to integer Euclidean division, both M and R are uniquely defined by these relations.

Notation and Symbols 2. Suppose that the polynomials A, B, M and R satisfy equation XX. We often use the following notation to describe the quotient and the remainder polynomials of the Euclidean division:

$$A \operatorname{div} B := M, \quad A \operatorname{mod} B := R \tag{3.26}$$

We also say that a polynomial A is divisible by another polynomial B if $A \mod B = 0$ holds. In this case, we also write B|A and call B a *factor* of A.

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Analogously to integers, methods to compute Euclidean division for polynomials are called **polynomial division algorithms**. Probably the best known algorithm is the so called **polynomial long division**.

Algorithm 3 Polynomial Euclidean Algorithm

```
Require: A, B \in R[x] with B \neq 0, such that Lc(B)^{-1} exists in R

procedure POLY-LONG-DIVISION(A, B)

M \leftarrow 0
R \leftarrow A
d \leftarrow deg(B)
c \leftarrow Lc(B)

while deg(R) \geq d do
S := Lc(R) \cdot c^{-1} \cdot x^{deg(R) - d}
M \leftarrow M + S
R \leftarrow R - S \cdot B
end while
return(Q, R)
end procedure

Ensure: A = M \cdot B + R
```

This algorithm works only when there is a notion of division by the leading coefficient of *B*. It can be generalized, but we will only need this somewhat simpler method in what follows.

Example 21 (Polynomial Long Division). To give an example of how the previous algorithm works, let us divide the integer polynomial $A(x) = x^5 + 2x^3 - 9 \in \mathbb{Z}[x]$ by the integer polynomial $B(x) = x^2 + 4x - 1 \in \mathbb{Z}[x]$. Since B is not the zero polynomial and the leading coefficient of B is 1, which is invertible as an integer, we can apply algorithm 1. Our goal is to find solutions to equation XXX, that is, we need to find the quotient polynomial $M \in \mathbb{Z}[x]$ and the reminder polynomial $R \in \mathbb{Z}[x]$ such that $x^5 + 2x^3 - 9 = M(x) \cdot (x^2 + 4x - 1) + R$. Using a notation that is mostly used in anglophone countries, we compute as follows:

$$\begin{array}{r}
X^{3} - 4X^{2} + 19X - 80 \\
X^{5} + 2X^{3} - 9 \\
\underline{-X^{5} - 4X^{4} + X^{3}} \\
-4X^{4} + 3X^{3} \\
\underline{4X^{4} + 16X^{3} - 4X^{2}} \\
\underline{-19X^{3} - 76X^{2} + 19X} - 9 \\
\underline{-80X^{2} + 19X - 9} \\
\underline{-80X^{2} + 320X - 80} \\
339X - 89
\end{array}$$
(3.27)

We therefore get $M(x) = x^3 - 4x^2 + 19x - 80$ as well as R(x) = 339x - 89 and indeed we have $x^5 + 2x^3 - 9 = (x^3 - 4x^2 + 19x - 80) \cdot (x^2 + 4x - 1) + (339x - 89)$, which we can double check invoking Sage:

```
1524 sage: Zx = ZZ['x']
1525 sage: A = Zx([-9,0,0,2,0,1]) 126
```

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 sage:
$$B = Zx([-1,4,1])$$
 127

 1527
 sage: $M = Zx([-80,19,-4,1])$
 128

 1528
 sage: $R = Zx([-89,339])$
 129

 1529
 sage: $A == M*B + R$
 130

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 True
 131

Example 22. In the previous example, polynomial division gave a non-trivial (non-vanishing, i.e non-zero) remainder. Of special interest are divisions that don't give a remainder. Such divisors are called factors of the dividend.

For example, consider the integer polynomial P_7 from example XXX again. As we have shown, it can be written both as $x^3 - 4x^2 - 11x + 30$ and as (x-2)(x+3)(x-5). From this, we can see that the polynomials $F_1(x) = (x-2)$, $F_2(x) = (x+3)$ and $F_3(x) = (x-5)$ are all factors of $x^3 - 4x^2 - 11x + 30$, since division of P_7 by any of these factors will result in a zero remainder.

Exercise 23. Consider the polynomial expressions $P(x) := -3x^4 + 4x^3 + 2x^2 + 4$ and $Q(x) = x^2 - 4x + 2$. Compute the Euclidean division of P by Q in the following types:

- 1. $P,Q \in \mathbb{Z}[x]$
- 1542 2. $P, Q \in \mathbb{Z}_6[x]$

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1543 3. $P, Q \in \mathbb{F}_5[x]$

Now consider the result in $\mathbb{Z}[x]$ and in $\mathbb{Z}_6[x]$. How can we compute the result in $\mathbb{Z}_6[x]$ from the result in $\mathbb{Z}[x]$?

Exercise 24. Show that the polynomial $P(x) = 2x^4 - 3x + 4 \in \mathbb{F}_5[x]$ is a factor of the polynomial $Q(x) = x^7 + 4x^6 + 4x^5 + x^3 + 2x^2 + 2x + 3 \in \mathbb{F}_5[x]$ that vis show P|Q. What is Q div P?

Prime Factors Recall that the fundamental theorem of arithmetics XXX tells us that every number is the product of prime numbers. Something similar holds for polynomials, too.

The polynomial analog to a prime number is a so called an **irreducible polynomial**, which is defined as a polynomial that cannot be factored into the product of two non-constant polynomials using Euclidean division. Irreducible polynomials are for polynomials what prime numbers are for integer: They are the basic building blocks from which all other polynomials can be constructed. To be more precise, let $P \in R[x]$ be any polynomial. Then there are always irreducible polynomials $F_1, F_2, \ldots, F_k \in R[x]$, such that the following holds:

$$P = F_1 \cdot F_2 \cdot \ldots \cdot F_k \ . \tag{3.28}$$

This representation is unique, except for permutations in the factors and is called the **prime** factorization of P.

Example 23. Consider the polynomial expression $P = x^2 - 3$. When we interpret P as an integer polynomial $P \in \mathbb{Z}[x]$, we find that this polynomial is irreducible, since any factorization other then $1 \cdot (x^2 - 3)$, must look like (x - a)(x + a) for some integer a, but there is no integers a with $a^2 = 3$.

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On the other hand interpreting P as a polynomial $P \in \mathbb{Z}_6[x]$ in modulo 6 arithmetics, we see that P has two factors $F_1 = (x-3)$ and $F_2 = (x+3)$, since $(x-3)(x+3) = x^2 - 3x + 3 - 3 \cdot 3 = x^2 - 3$.

Finding prime factors of a polynomial is hard. As we have seen in example XXX, points where a polynomial evaluates to zero, i.e points $x_0 \in R$ with $P(x_0) = 0$ are of special interest, since it can be shown the polynomial $F(x) = (x - x_0)$ is always a factor of P. The converse, however, is not necessarily true, because a polynomial can have irreducible prime factors.

Points where a polynomial evaluates to zero are called the **roots** of the polynomial. To be more precise, let $P \in R[x]$ be a polynomial. Then the set of all roots of P is defined as follows:

$$R_0(P) := \{ x_0 \in R \mid P(x_0) = 0 \}$$
(3.29)

Finding the roots of a polynomial is sometimes called **solving the polynomial**. It is a hard problem and has been the subject of much research throughout history. In fact, it is well known that, for polynomials of degree 5 or higher, there is, in general, no closed expression, from which the roots can be deduced.

It can be shown that if m is the degree of a polynomial P, then P can not have more than m roots. However, in general, polynomials can have less than m roots.

Example 24. Consider our integer polynomial $P_7(x) = x^3 - 4x^2 - 11x + 30$ from example XXX again. We know that its set of roots is given by $R_0(P_7) = \{-3, 2, 5\}$.

On the other hand, we know from example XXX that the integer polynomial $x^2 - 3$ is irreducible. It follows that it has no roots, since every root defines a prime factor.

Example 25. To give another example, consider the integer polynomial $P = x^7 + 3x^6 + 3x^5 + x^4 - x^3 - 3x^2 - 3x - 1$. We can invoke Sage to compute the roots and prime factors of P:

We see that P has the root 1 and that the associated prime factor (x-1) occurs once in P and that it has the root -1, where the associated prime factor (x+1) occurs 4 times in P. This gives the prime following factorization:

$$P = (x-1)(x+1)^4(x^2+1)$$

Lange interpolation One particularly usefu property of polynomials is that a polynomial of degree m is completely determined on m+1 evaluation points. Seeing this from a different angle, we can (sometimes) uniquely derive a polynomial of degree m from a set S:

$$S = \{(x_0, y_0), (x_1, y_1), \dots, (x_m, y_m) \mid x_i \neq x_j \text{ for all indices i and j} \}$$
 (3.30)

This "few too many" property of polynomials is used in many places, like for example in erasure codes. It is also of importance in snarks and we therefore need to understand a method to actually compute a polynomial from a set of points.

what does this mean?

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If the coefficients of the polynomial we want to find have a notion of multiplicative inverse, it is always possible to find such a polynomial. One method for this is called **Lagrange interpolation**. It works as follows: Given a set like 3.30, a polynomial P of degree m+1 with $P(x_i) = y_i$ for all pairs (x_i, y_i) from S is given by the following algorithm:

```
Algorithm 4 Lagrange Interpolation

Require: R must have multiplicative inverses

Require: S = \{(x_0, y_0), (x_1, y_1), \dots, (x_m, y_m) \mid x_i, y_i \in R, x_i \neq x_j \text{ for all indices i and j} \}

procedure LAGRANGE-INTERPOLATION(S)

for j \in (0 \dots m) do

l_j(x) \leftarrow \prod_{i=0; i \neq j}^m \frac{x-x_i}{x_j-x_i} = \frac{(x-x_0)}{(x_j-x_0)} \cdots \frac{(x-x_{j-1})}{(x_j-x_{j-1})} \frac{(x-x_{j+1})}{(x_j-x_{j+1})} \cdots \frac{(x-x_m)}{(x_j-x_m)}

end for
P \leftarrow \sum_{j=0}^m y_j \cdot l_j
return P
end procedure
```

Ensure: $P \in R[x]$ with deg(P) = m

Ensure: $P(x_j) = y_j$ for all pairs $(x_j, y_j) \in S$

Example 26. Let us consider the set $S = \{(0,4), (-2,1), (2,3)\}$. Our task is to compute a polynomial of degree 2 in $\mathbb{Q}[x]$ with fractional number coefficients. Since \mathbb{Q} has multiplicative inverses, we can use the Lagrange interpolation algorithm from 4, to compute the polynomial.

$$l_0(x) = \frac{x - x_1}{x_0 - x_1} \cdot \frac{x - x_2}{x_0 - x_2} = \frac{x + 2}{0 + 2} \cdot \frac{x - 2}{0 - 2} = -\frac{(x + 2)(x - 2)}{4}$$

$$= -\frac{1}{4}(x^2 - 4)$$

$$l_1(x) = \frac{x - x_0}{x_1 - x_0} \cdot \frac{x - x_2}{x_1 - x_2} = \frac{x - 0}{-2 - 0} \cdot \frac{x - 2}{-2 - 2} = \frac{x(x - 2)}{8}$$

$$= \frac{1}{8}(x^2 - 2x)$$

$$l_2(x) = \frac{x - x_0}{x_2 - x_0} \cdot \frac{x - x_1}{x_2 - x_1} = \frac{x - 0}{2 - 0} \cdot \frac{x + 2}{2 + 2} = \frac{x(x + 2)}{8}$$

$$= \frac{1}{8}(x^2 + 2x)$$

$$P(x) = 4 \cdot (-\frac{1}{4}(x^2 - 4)) + 1 \cdot \frac{1}{8}(x^2 - 2x) + 3 \cdot \frac{1}{8}(x^2 + 2x)$$

$$= -x^2 + 4 + \frac{1}{8}x^2 - \frac{1}{4}x + \frac{3}{8}x^2 + \frac{3}{4}x$$

$$= -\frac{1}{2}x^2 + \frac{1}{2}x + 4$$

And, indeed, evaluation of P on the x-values of S gives the correct points, since P(0) = 4, P(-2) = 1 and P(2) = 3.

Example 27. To give another example more relevant to the topics of this book, let us consider the same set $S = \{(0,4), (-2,1), (2,3)\}$ as in the previous example. This time, the task is to compute a polynomial $P \in \mathbb{F}_5[x]$ from this data. Since we know that multiplicative inverses exist in \mathbb{F}_5 , algorithm XXX applies and we can compute a unique polynomial of degree 2 in $\mathbb{F}_5[x]$

from S. We can use the lookup tables XXX for computation in \mathbb{F}_5 and get the following:

$$l_0(x) = \frac{x - x_1}{x_0 - x_1} \cdot \frac{x - x_2}{x_0 - x_2} = \frac{x + 2}{0 + 2} \cdot \frac{x - 2}{0 - 2} = \frac{(x + 2)(x - 2)}{-4} = \frac{(x + 2)(x + 3)}{1}$$

$$= x^2 + 1$$

$$l_1(x) = \frac{x - x_0}{x_1 - x_0} \cdot \frac{x - x_2}{x_1 - x_2} = \frac{x - 0}{-2 - 0} \cdot \frac{x - 2}{-2 - 2} = \frac{x}{3} \cdot \frac{x + 3}{1} = 2(x^2 + 3x)$$

$$= 2x^2 + x$$

$$l_2(x) = \frac{x - x_0}{x_2 - x_0} \cdot \frac{x - x_1}{x_2 - x_1} = \frac{x - 0}{2 - 0} \cdot \frac{x + 2}{2 + 2} = \frac{x(x + 2)}{3} = 2(x^2 + 2x)$$

$$= 2x^2 + 4x$$

$$P(x) = 4 \cdot (x^2 + 1) + 1 \cdot (2x^2 + x) + 3 \cdot (2x^2 + 4x)$$

$$= 4x^2 + 4 + 2x^2 + x + x^2 + 2x$$

$$= 2x^2 + 3x + 4$$

And, indeed, evaluation of P on the x-values of S gives the correct points, since P(0)=4, P(-2)=1 and P(2)=3.

Exercise 25. Consider example XXX and example XXX again. Why is it not possible to apply algorithm XXX if we consider S as a set of integers, nor as a set in \mathbb{Z}_6 ?

Chapter 4

Algebra

In the previous chapter, we gave an introduction to the basic computational skills needed for a pen-and-paper approach to SNARKs. This chapter provides a more abstract clarification of relevant mathematical terminology on **algebraic types** such as **groups**, **fields**, **rings** and similar.

In a nutshell, algebraic types define sets that are analogous to numbers in various aspects, in the sense that you can add, subtract, multiply or divide on those sets. We know many examples of sets that fall under those categories, such as natural numbers, integers, rational or the real numbers. In some sense, these are the most fundamental examples of such sets.

Papers on cryptography (and mathematical papers in general) frequently contain such terms, and it is necessary to get at least some understanding of these terms to be able to follow these papers. In this chapter, we therefore provide a short introduction to these concepts.

Def Subgroup, Fundamental theorem of cyclic groups.

4.1 Groups

Groups are abstractions that capture the essence of mathematical phenomena, like addition and subtraction, multiplication and division, permutations, or symmetries.

To understand groups, let us think back to when we learned about the addition and subtraction of integers (also called whole numbers) in school. We have learned that, whenever we add two integers, the result is guaranteed to be an integer as well. We have also learned that adding zero to any integer means that "nothing happens", that is, result of the addition is the same integer we started with. Furthermore, we have learned that the order in which we add two (or more) integers does not matter, that operations within brackets should be computed before operations outside brackets, and that, for every integer, there is always another integer (the negative) such that we get zero when we add them together.

These conditions are the defining properties of a group, and mathematicians have recognized that the exact same set of rules can be found in very different mathematical structures. It therefore makes sense to give a formal definition of what a group should be, detached from any concrete examples. This lets us handle entities of very different mathematical origins in a flexible way, while retaining essential structural aspects of many objects in abstract algebra and beyond.

Distilling these rules to the smallest independent list of properties and making them abstract, we arrive at the definition of a group:

Definition 4.1.0.1. A **group** (\mathbb{G},\cdot) is a set \mathbb{G} , together with a **map** \cdot . The map, also denoted as $\mathbb{G} \times \mathbb{G} \to \mathbb{G}$ and called the **group law**, combines two elements of the set \mathbb{G} into a third one such that the following properties hold:

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- Existence of a neutral element: There is a $e \in \mathbb{G}$ for all $g \in \mathbb{G}$, such that $e \cdot g = g$ as well as $g \cdot e = g$.
- Existence of an inverse: For every $g \in \mathbb{G}$ there is a $g^{-1} \in \mathbb{G}$, such that $g \cdot g^{-1} = e$ as well as $g^{-1} \cdot g = e$.
- Associativity: For every $g_1, g_2, g_3 \in \mathbb{G}$ the equation $g_1 \cdot (g_2 \cdot g_3) = (g_1 \cdot g_2) \cdot g_3$ holds.

Rephrasing the abstract definition in layman's terms, a group is something where we can do computations in a way that resembles the behavior of the addition of integers. Specifically, this means we can combine some element with another element into a new element in a way that is reversible and where the order of combining elements doesn't matter.

Notation and Symbols 3. Let $(\mathbb{G}\cdot)$ be a finite group. If there is no risk of ambiguity (whether we are talking about agroup or a set), we frequently drop the symbol \cdot and simply write \mathbb{G} as the notation for the group, keeping the group law implicit.

As we will see in XXX, groups are heavily used in cryptography and in SNARKs. But let us look at some more familiar examples fist:

add reference

Example 28 (Integer Addition and Subtraction). The set $(\mathbb{Z},+)$ of integers with integer addition is the archetypical example of a group, where the group law is traditionally written as + (instead of \cdot). To compare integer addition against the abstract axioms of a group, we first see that the neutral element e is the number 0, since a+0=a for all integers $a \in \mathbb{Z}$. Furthermore, the inverse of a number is its negative counterpart, since a+(-a)=0, for all $a \in \mathbb{Z}$. In addition, we know that (a+b)+c=a+(b+c), so integers with addition are indeed a group in the abstract sense.

Example 29 (The trivial group). The most basic example of a group is group with just one element $\{\bullet\}$ and the group law $\bullet \cdot \bullet = \bullet$.

Add reallife example of

Commutative Groups When we look at the general definition of a group, we see that it is somewhat different from what we know from integers. We know that the order in which we add two integers doesn't matter, as, for example, 4+2 is the same as 2+4. However, we also know from example XXX that this is not the case for all groups.

add reference

This means that groups where the order in which the group law is executed doesn't matter are a special subcase of groups called **commutative groups**. To be more precise, a group is called commutative if $g_1 \cdot g_2 = g_2 \cdot g_1$ holds for all $g_1, g_2 \in \mathbb{G}$.

Notation and Symbols 4. For commutative groups (\mathbb{G},\cdot) , we frequently use the so-called **addi**tive notation $(\mathbb{G},+)$, that is, we write + instead of \cdot for the group law, and $-g:=g^{-1}$ for the inverse of an element $g\in\mathbb{G}$.

Example 30. Consider the group of integers with integer addition again. Since a + b = b + a for all integers, this group is the archetypical example of a commutative group. Since there are infinitely many integers, $(\mathbb{Z}, +)$ is not a finite group.

Example 31. Consider our definition of modulo 6 residue classes $(\mathbb{Z}_6,+)$ as defined in the addition table from example 8. As we can see, the residue class 0 is the neutral element in modulo 6 arithmetics, and the inverse of a residue class r is given by 6-r, since r+(6-r)=6, which is congruent to 0, since 6 mod 6=0. Moreover, $(r_1+r_2)+r_3=r_1+(r_2+r_3)$ is inherited from integer arithmetic.

We therefore see that $(\mathbb{Z}_6, +)$ is a group, and, since the addition table in example 8 is symmetrical, we see $r_1 + r_2 = r_2 + r_1$, which shows that $(\mathbb{Z}_6, +)$ is commutative.

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The previous example of a commutative group is a very important one for this book. Abstracting from this example and considering residue classes $(\mathbb{Z}_n, +)$ for arbitrary moduli n, it can be shown that $(\mathbb{Z}, +)$ is a commutative group with the neutral element 0 and the additive inverse n-r for any element $r \in \mathbb{Z}_n$. We call such a group the **remainder class group** of modulus n.

Of particular importance for pairing-based cryptography in general and SNARKs in particular are so-called **pairing maps** on commutative groups. To be more precise let \mathbb{G}_1 , \mathbb{G}_2 and \mathbb{G}_3 be three commutative groups. For historical reasons, we write the group law on \mathbb{G}_1 and \mathbb{G}_2 in additive notation and the group law on \mathbb{G}_3 in multiplicative notation. Then a **pairing map** is a function

$$e(\cdot,\cdot): \mathbb{G}_1 \times \mathbb{G}_2 \to \mathbb{G}_3$$
 (4.1)

This function takes pairs (g_1, g_2) (products) of elements from \mathbb{G}_1 and \mathbb{G}_2 , and maps them to elements from \mathbb{G}_3 , such that the **bilinearity** property holds:

1700 Definition 4.1.0.2. Bilinearity

For all $g_1, g_1' \in \mathbb{G}_1$ and $g_2 \in \mathbb{G}_2$ we have $e(g_1 + g_1', g_2) = e(g_1, g_2) \cdot e(g_1', g_2)$ and for all $g_1 \in \mathbb{G}_1$ and $g_2, g_2' \in \mathbb{G}_2$ we have $e(g_1, g_2 + g_2') = e(g_1, g_2) \cdot e(g_1, g_2')$.

A pairing map is called **non-degenerate** if, whenever the result of the pairing is the neutral element in \mathbb{G}_3 , one of the input values is the neutral element of \mathbb{G}_1 or \mathbb{G}_2 . To be more precise, $e(g_1,g_2)=e_{\mathbb{G}_3}$ implies $g_1=e_{\mathbb{G}_1}$ or $g_2=e_{\mathbb{G}_2}$.

Informally speaking, bilinearity means that it doesn't matter if we first execute the group law on one side and then apply the bilinear map, or if we first apply the bilinear map and then apply the group law. Moreover, non-degeneracy means that the result of the pairing is zero if and only if at least one of the input values is zero.

Example 32. One of the most basic examples of a non-degenerate pairing involves \mathbb{G}_1 , \mathbb{G}_2 and \mathbb{G}_3 all to be the groups of integers with addition $(\mathbb{Z},+)$. Then the following map defines a non-degenerate pairing:

$$e(\cdot,\cdot): \mathbb{Z} \times \mathbb{Z} \to \mathbb{Z} \ (a,b) \mapsto a \cdot b$$

Note that bilinearity follows from the distributive law of integers, since for $a,b,c\in\mathbb{Z}$, we have $e(a+b,c)=(a+b)\cdot c=a\cdot c+b\cdot c=e(a,c)+e(b,c)$ and the same reasoning is true for the second argument.

To the see that $e(\cdot, \cdot)$ is non-degenrate, assume that e(a, b) = 0. Then $a \cdot b = 0$ implies that a or b must be zero.

Exercise 26. Consider example 13 again and let \mathbb{F}_5^* be the set of all remainder classes from \mathbb{F}_5 without the class 0. Then $\mathbb{F}_5^* = \{1, 2, 3, 4\}$. Show that (\mathbb{F}_5^*, \cdot) is a commutative group.

check reference

Exercise 27. Generalizing the previous exercise, consider the general modulus n, and let \mathbb{Z}_n^* be the set of all remainder classes from \mathbb{Z}_n without the class 0. Then $\mathbb{Z}_n^* = \{1, 2, ..., n-1\}$. Provide a counter-example to show that (\mathbb{Z}_n^*, \cdot) is not a group in general.

Find a condition such that (\mathbb{Z}_n^*, \cdot) is a commutative group, compute the neutral element, give a closed form for the inverse of any element and prove the commutative group axioms.

Exercise 28. Consider the remainder class groups $(\mathbb{Z}_n,+)$ for some modulus n. Show that the map

$$e(\cdot,\cdot): \mathbb{Z}_n \times \mathbb{Z}_n \to \mathbb{Z}_n \ (a,b) \mapsto a \cdot b$$

is bilinear. Why is it not a pairing in general and what condition must be imposed on n, such that the map will be a pairing?

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Finite groups As we have seen in the previous examples, groups can either contain infinitely many elements (such as integers) or finitely many elements (as for example the remainder class groups $(\mathbb{Z}_n, +)$). To capture this distinction, a group is called a **finite group** if the underlying set of elements is finite. In that case, the number of elements of that group is called its **order**.

Notation and Symbols 5. Let \mathbb{G} be a finite group. We write $ord(\mathbb{G})$ or $|\mathbb{G}|$ for the order of \mathbb{G} .

Example 33. Consider the remainder class groups $(\mathbb{Z}_6,+)$ from example 8 and $(\mathbb{F}_5,+)$ from example 13, and the group (\mathbb{F}_5^*,\cdot) from exercise 26. We can easily see that the order of $(\mathbb{Z}_6,+)$ is 6, the order of $(\mathbb{F}_5,+)$ is 5 and the order of (\mathbb{F}_5^*,\cdot) is 4.

To be more general, considering arbitrary moduli n, we know from Euclidean division that the order of the remainder class group $(\mathbb{Z}_n, +)$ is n.

Exercise 29. The RSA crypto system is based on a modulus n that is typically the product of two prime numbers of size 2048-bits. What is the approximate order of the remainder class group $(\mathbb{Z}_n, +)$ in this case?

Generators These are sets of elements that can be used to generate the entire group by applying the group law repeatedly to these elements or their inverses only. Generators are of particular interest when working with groups.

Of course, every group \mathbb{G} has a trivial set of generators, when we just consider every element of the group to be in the generator set. The more interesting question is to find the smallest possible set of generators for a given group. Of particular interest in this regard are groups that have a single generator, that is, there exists an element $g \in \mathbb{G}$ such that every other element from \mathbb{G} can be computed by the repeated combination of g and its inverse g^{-1} only. Groups with a single generator are called **cyclic groups**.

Example 34. The most basic example of a cyclic group is the group of integers $(\mathbb{Z},+)$ with integer addition. 1 is a single generator of \mathbb{Z} , since every integer can be obtained by repeatedly adding either 1 or its inverse -1 to itself. For example -4 is generated by -1, since -4 = -1 + (-1) + (-1) + (-1).

Example 35. Consider a modulus n and the remainder class groups $(\mathbb{Z}_n, +)$ from example 33. These groups are cyclic, with the generator 1, since every other element of that group can be constructed by repeatedly adding the remainder class 1 to itself. Since \mathbb{Z}_n is also finite, we know that $(\mathbb{Z}_n, +)$ is a finite cyclic group of order n.

Example 36. Let $p \in \mathbb{P}$ be prime number and (\mathbb{F}_p^*, \cdot) the finite group from exercise XXX. Then (\mathbb{F}_p^*, \cdot) is cyclic and every element $g \in \mathbb{F}_q^*$ is a generator.

The discrete Logarithm problem Observe that, when \mathbb{G} is a cyclic group of order n and $g \in \mathbb{G}$ is a generator of \mathbb{G} , then there is a map with respect to the generator g with the following properties:

$$g^{(\cdot)}: \mathbb{Z}_n \to \mathbb{G} \ x \mapsto g^x$$
 (4.2)

In the map above, g^x means "multiply g x-times by itself" and $g^0 = e_{\mathbb{G}}$. This map, called the **exponential map**, has the remarkable property that it maps the additive group law of the remainder class group $(\mathbb{Z}_n, +)$ in a one-to-one correspondence to the group law of \mathbb{G} .

To see this, first observe that, since $g^0 := e_{\mathbb{G}}$ by definition, the neutral element of \mathbb{Z}_n is mapped to the neutral element of \mathbb{G} , and, since $g^{x+y} = g^x \cdot g^y$, the map respects the group law.

Because the exponential map respects the group law, it doesn't matter if we do our computation in \mathbb{Z}_n before we write the result into the exponent of g or afterwards: the result will be the

check references to previous examples

RSA crypto system

size 2048-bits

add reference: 27?

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reference

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same in both cases. This is usually referred to as doing computations "in the exponent". In cryptography in general, and in SNARK development in particular, we often perform computations "in the exponent" of a generator.

Example 37. Consider the multiplicative group (\mathbb{F}_5^*, \cdot) from example 26. We know that \mathbb{F}_5^* is a cyclic group of order 4, and that every element is a generator. If we choose $3 \in \mathbb{F}_5^*$, we then know that the following map respects the group law of addition in \mathbb{Z}_4 and the group law of multiplication in \mathbb{F}_5^* :

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$$3^{(\cdot)}: \mathbb{Z}_4 \to \mathbb{F}_5^* x \mapsto 3^x$$

Let us now perform a computation in the exponent:

$$3^{2+3-2} = 3^3$$

= 2

This gives the same result as doing the same computation in $\mathbb{F}*_5$:

$$3^{2+3-2} = 3^{2} \cdot 3^{3} \cdot 3^{-2}$$

$$= 4 \cdot 2 \cdot (-3)^{2}$$

$$= 3 \cdot 2^{2}$$

$$= 3 \cdot 4$$

$$= 2$$

Since the exponential map is a one-to-one correspondence that respects the group law, it can be shown that this map has an inverse with respect to the base g, called the **discrete logarithm** map:

$$log_g(\cdot): \mathbb{G} \to \mathbb{Z}_n \, x \mapsto log_g(x)$$
 (4.3)

Discrete logarithms are highly important in cryptography, because there are groups such that the exponential map and its inverse, the discrete logarithm, which are believed to be one-way functions, that is, while it is possible to compute the exponential map in polynomial time, computing the discrete log takes (sub)-exponential time. We have discussed this briefly following example 3.5 in the previous chapter, and will look at this and similar problems in more detail in the next section.

polynomial time

exponentia time

4.1.1 Cryptographic Groups

In this section, we will look at families of groups that are believed to satisfy certain **computational hardness assumptions**, namely that a particular problem cannot be solved efficiently (where efficiently typically means "in polynomial time of a given security parameter") in the groups under consideration.

Example 38. To highlight the concept of the computational hardness assumption, consider the group of integers \mathbb{Z} from example 3.5. One of the best known and most researched examples of computational hardness is the assumption that the factorization of integers into prime numbers cannot be solved by any algorithm in polynomial time with respect to the bit-length of the integer.

To be more precise, the computational hardness assumption of integer factorization assumes that, given any integer $z \in \mathbb{Z}$ with bit-length b, there is no integer k and no algorithm with the

mental theorem of finite cyclic groups

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runtime complexity of $\mathcal{O}(b^k)$ that is able to find the prime numbers $p_1, p_2, \dots, p_j \in \mathbb{P}$, such that $z = p_1 \cdot p_2 \cdot \dots \cdot p_j$.

This hardness assumption was proven to be false, since Shor's (1994) algorithm shows that integer factorization is at least efficiently possible on a quantum computer, since the runtime complexity of this algorithm is $\mathcal{O}(b^3)$. However, no such algorithm is known on a classical computer.

In the realm of classical computers, however, we still have to call the non-existence of such an algorithm an "assumption" because, to date, there is no proof that it is actually impossible to find one. The problem is that it is hard to reason about algorithms that we don't know.

So, despite the fact that there is currently no known algorithm that can factor integers efficiently on a classical computer, we cannot exclude that such an algorithm might exist in principle, and that someone eventually will discover it in the future.

However, what still makes the assumption plausible, despite the absence of any actual proof, is the fact that, after decades of extensive search, still no such algorithm has been found.

In what follows, we will describe a few <u>computational hardness assumptions that arise in</u> the context of groups in cryptography, because we will refer to them throughout the book.

The discrete logarithm assumption The so-called discrete logarithm problem is one of the most fundamental assumptions in cryptography. To define it, let \mathbb{G} be a finite cyclic group of order r and let g be a generator of \mathbb{G} . We know from 4.2 that there is an exponential map $g^{(\cdot)}: \mathbb{Z}_r \to \mathbb{G}: x \mapsto g^x$ that maps the residue classes from modulo r arithmetic onto the group in a 1:1 correspondence. The **discrete logarithm problem** is the task of finding inverses to this map, that is, to find a solution $x \in \mathbb{Z}_r$ to the following equation for some given $h \in \mathbb{G}$:

$$h = g^{x} \tag{4.4}$$

In other words, the **discrete logarithm assumption** (**DL-A**) is the assumption that there exists no algorithm with polynomial running time in the security parameter $log_2(r)$, that is able to compute some x if only h, g and g^x are given in \mathbb{G} . If this is the case for \mathbb{G} , we call \mathbb{G} a **DL-A group**.

Rephrasing the previous definition, DL-A groups are believed to have the property that it is infeasible to compute some number x that solves the equation $h = g^x$ for a given h and g, assuming that the size of the group r is large enough.

Example 39 (Public key cryptography). One the most basic examples of an application for DL-A groups is in public key cryptography, where the parties publicly agree on some pair (\mathbb{G}, g) such that \mathbb{G} is a finite cyclic group of sufficiently large order r, where \mathbb{G} is believed to be a DL-A group, and g is a generator of \mathbb{G} .

In this setting, a secret key is some number $sk \in \mathbb{Z}_r$ and the associated public key pk is the group element $pk = g^{sk}$. Since discrete logarithms are assumed to be hard, it is infeasible for an attacker to compute the secret key from the public key, since it is believed to be hard to find solutions x to the following equation:

$$pk = g^{x} (4.5)$$

As the previous example shows, identifying DL-A groups is an important practical problem. Unfortunately, it is easy to see that it does not make sense to assume the hardness of the discrete logarithm problem in all finite cyclic groups: Counterexamples are common and easy to construct.

runtime complexity

S: what does "efficiently" mean here?

hardness assumptions

computationa

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Example 40 (Modular arithmetics for Fermat's primes). It is widely believed that the discrete logarithm problem is hard in multiplicative groups \mathbb{Z}_p^* of prime number modular arithmetics. However, this is not true in general. To see that, consider any so-called Fermat's prime, which is a prime number $p \in \mathbb{P}$, such that $p = 2^n + 1$ for some number n.

We know from exercise 27 that in this case $\mathbb{Z}_p^* = \{1, 2, ..., p-1\}$ is a group with respect to integer multiplication in modular p arithmetics and since $p = 2^n + 1$, the order of \mathbb{Z}_p^* is 2^n , reference to integer multiplication in modular p arithmetics and since $p = 2^n + 1$, the order of \mathbb{Z}_p^* is 2^n , which implies that the associated security parameter is given by $log_2(2^n) = n$.

We show that, in this case, \mathbb{Z}_p^* is not a DL-A group, by constructing an algorithm, which is able compute some $x \in \mathbb{Z}_{2^n}$ for any given generator g and arbitrary element h of \mathbb{F}_p^* , such that equation 4.6 holds, and the runtime complexity of the constructed algorithm is $\mathcal{O}(n^2)$, which is quadratic in the security parameter $n = log_2(2^n)$.

$$h = g^{x} \tag{4.6}$$

To define such an algorithm, let us assume that the generator g is a public constant and that a group element h is given. Our task is to compute x efficiently.

The first thing to note is that, since x is a number in modular 2^n arithmetic, we can write the binary representation of x as in 4.7, with binary coefficients $c_i \in \{0,1\}$. In particular, x is an *n*-bit number if interpreted as an integer.

$$x = c_0 \cdot 2^0 + c_1 \cdot 2^1 + \dots + c_n \cdot 2^n$$
(4.7)

explain last sentence more

We then use this representation to construct an algorithm that computes the bits c_i one after another, starting at c_0 . To see how this can be achieved, observe that we can determine c_0 by raising the input h to the power of 2^{n-1} in \mathbb{F}_p^* . We use the exponential laws and compute as follows:

$$h^{2^{n-1}} = (g^{x})^{2^{n-1}}$$

$$= (g^{c_0 \cdot 2^0 + c_1 \cdot 2^1 + \dots + c_n \cdot 2^n})^{2^{n-1}}$$

$$= g^{c_0 \cdot 2^{n-1}} \cdot g^{c_1 \cdot 2^1 \cdot 2^{n-1}} \cdot g^{c_2 \cdot 2^2 \cdot 2^{n-1}} \cdots g^{c_n \cdot 2^n \cdot 2^{n-1}}$$

$$= g^{c_0 \cdot 2^{n-1}} \cdot g^{c_1 \cdot 2^0 \cdot 2^n} \cdot g^{c_2 \cdot 2^1 \cdot 2^n} \cdots g^{c_n \cdot 2^{n-1} \cdot 2^n}$$

Now, since g is a generator and \mathbb{F}_p^* is cyclic of order 2^n , we know $g^{2^n} = 1$ and therefore $g^{k \cdot 2^n} = 1$ 1849 $1^k = 1$. From this, it follows that all but the first factor in the last expression are equal to 1 and 1850 we can simplify the expression into the following: 1851

$$h^{2^{n-1}} = g^{c_0 2^{n-1}} (4.8)$$

Now, in case $c_0 = 0$, we get $h^{2^{n-1}} = g^0 = 1$. In case $c_0 = 1$, we get $h^{2^{n-1}} = g^{2^{n-1}} \neq 1$ (To see that $g^{2^{n-1}} \neq 1$, recall that g is a generator of \mathbb{F}_p^* and hence, is \mathbb{F}_p^* a cyclic group of order 2^n , which implies $g^y \neq 1$ for all $y < 2^n$).

Raising h to the power of 2^{n-1} determines c_0 , and we can apply the same reasoning to the coefficient c_1 by raising $h \cdot g^{-c_0 \cdot 2^0}$ to the power of 2^{n-2} . This approach can then be repeated until all the coefficients c_i of x are found.

Assuming that exponentiation in \mathbb{F}_p^* can be done in logarithmic runtime complexity log(p), it follows that our algorithm has a runtime complexity of $\mathcal{O}(\log^2(p)) = \mathcal{O}(n^2)$, since we have to execute n exponentiations to determine the n binary coefficients of x.

From this, it follows that whenever p is a Fermat's prime, the discrete logarithm assumption does not hold in F_p^* .

The decisional Diffie–Hellman assumption Let \mathbb{G} be a finite cyclic group of order r and let g be a generator of \mathbb{G} . The decisional Diffie–Hellman assumption stipulates that there is no algorithm that has a polynomial runtime complexity in the security parameter s = log(r) that is able to distinguish the so-called DDH- triple (g^a, g^b, g^{ab}) from any triple (g^a, g^b, g^c) for randomly and independently chosen parameters $a, b, c \in \mathbb{Z}_r$. If this is the case for \mathbb{G} , we call \mathbb{G} a **DDH-A group**.

DDH-A is a stronger assumption than DL-A, in the sense that the discrete logarithm assumption is necessary for the decisional Diffie–Hellman assumption to hold, but not the other way around.

To see why this is the case, assume that the discrete logarithm assumption does not hold. In that case, given a generator g and a group element h, it is easy to compute some residue class $x \in \mathbb{Z}_p$ with $h = g^x$. Then the decisional Diffie-Hellman assumption cannot hold, since given some triple (g^a, g^b, z) , one could efficiently decide whether $z = g^{ab}$ is true by first computing the discrete logarithm b of g^b , then computing $g^{ab} = (g^a)^b$ and deciding whether or not $z = g^{ab}$.

On the other hand, the following example shows that there are groups where the discrete logarithm assumption holds but the decisional Diffie–Hellman assumption does not.

Example 41 (Efficiently computable pairings). Let \mathbb{G} be a finite, cyclic group of order r with generator g, such that the discrete logarithm assumtion holds and there is a pairing map $e(\cdot, \cdot)$: $\mathbb{G} \times \mathbb{G} \to \mathbb{G}_T$ for some target group \mathbb{G}_T that is computable in polynomial time of the parameter log(r).

In a setting like this, it is easy to show that DDH-A cannot hold, since given some triple (g^a, g^b, z) , it is possible to decide in polynomial time w.r.t log(r) whether $z = g^{ab}$ or not. To see that, check the following equation:

$$e(g^a, g^b) = e(g, z) \tag{4.9}$$

Since the bilinearity properties of $e(\cdot, \cdot)$ imply $e(g^a, g^b) = e(g, g)^{ab} = e(g, g^{ab})$, and e(g, y) = e(g, y') implies y = y' due to the non-degenerate property, the equality means $z = e^{ab}$.

"equation"

It follows that DDH-A is indeed weaker than DL-A, and groups with efficient pairings cannot be DDH-A groups. The following example shows another important class of groups where DDH-A does not hold: multiplicative groups of prime number residue classes.

Example 42. Let p be a prime number and $\mathbb{Z}_p^* = \{1, 2, ..., p-1\}$ the multiplicative group of modular p arithmetics as in exercise 27. As we have seen in XXX, this group is finite and cyclic of order p-1 and every element $g \neq 1$ is a generator.

To see that \mathbb{F}_p^* cannot be a DDH-A group, recall from XXX that the Legendre symbol $\left(\frac{x}{p}\right)$ of any $x \in \mathbb{F}_p^*$ is efficiently computable by Euler's formular. But the Legendre symbol of g^a reveals whether a is even or odd. Given g^a , g^b and g^{ab} , one can thus efficiently compute and compare the least significant bit of a, b and ab, respectively, which provides a probabilistic method to distinguish g^{ab} from a random group element g^c .

The computational Diffie–Hellman assumption Let \mathbb{G} be a finite cyclic group of order r and let g be a generator of \mathbb{G} . The computational Diffie–Hellman assumption stipulates that, given randomly and independently chosen residue classes $a, b \in \mathbb{Z}_r$, it is not possible to compute g^{ab} if only g, g^a and g^b (but not a and b) are known. If this is the case for \mathbb{G} , we call \mathbb{G} a CDH-A group.

In general, we don't know if CDH-A is a stronger assumption then DL-A, or if both assumptions are equivalent. We know that DL-A is necessary for CDH-A, but the other direction

check reference

what's the difference between \mathbb{F}_p^* and \mathbb{Z}_p^* ?

Legendre symbol

Euler's formular

These are only explained later in the text, '4.27'

is currently not well understood. In particular, there are no groups known where DL-A holds but CDH-A does not hold [Fifield, 2012].

To see why the discrete logarithm assumption is necessary, assume that it does not hold. So, given a generator g and a group element h, it is easy to compute some residue class $x \in \mathbb{Z}_p$ with $h = g^x$. In that case, the computational Diffie–Hellman assumption cannot hold, since, given g, g^a and g^b , one can efficiently compute b and hence is able to compute $g^{ab} = (g^a)^b$ from this data

The computational Diffie–Hellman assumption is a weaker assumption than the decisional Diffie–Hellman assumption, which means that there are groups where CDH-A holds and DDH-A does not hold, while there cannot be groups such that DDH-A holds but CDH-A does not hold. To see that, assume that it is efficiently possible to compute g^{ab} from g, g^a and g^b . Then, given (g^a, g^b, z) it is easy to decide if $z = g^{ab}$ holds or not.

Several variations and special cases of the CDH-A exist. For example, the **square computational Diffie–Hellman assumption** assumes that, given g and g^x , it is computationally hard to compute g^{x^2} . The **inverse computational Diffie–Hellman assumption** assumes that, given g and g^x , it is computationally hard to compute $g^{x^{-1}}$.

Cofactor Clearing

4.1.2 Hashing to Groups

Hash functions Generally speaking, a hash function is any function that can be used to map data of arbitrary size to fixed-size values. Since binary strings of arbitrary length are a general way to represent arbitrary data, we can understand a general **hash function** as the following map where $\{0,1\}^*$ represents the set of all binary strings of arbitrary but finite length and $\{0,1\}^k$ represents the set of all binary strings that have a length of exactly k bits:

$$H: \{0,1\}^* \to \{0,1\}^k$$
 (4.10)

In our definition, a hash function maps binary strings of arbitrary size onto binary strings of size exactly k. The **images** of H, that is, the values returned by the hash function H, are called **hash values**, **digests**, or simply **hashes**.

A hash function must be deterministic, that is, when we insert the same input x into H, the image H(x) must always be the same. In addition, a hash function should be as uniform as possible, which means that it should map input values as evenly as possible over its output range. In mathematical terms, every string of length k from $\{0,1\}^k$ should be generated with roughly the same probability.

Example 43 (k-truncation hash). One of the most basic hash functions $H_k : \{0,1\}^* \to \{0,1\}^k$ is given by simply truncating every binary string s of size s.len() > k to a string of size k and by filling any string s' of size s'.len() < 0 with zeros. To make this hash function deterministic, we define that both truncation and filling should happen "on the left".

For example, if k = 3, $x_1 = (000010101110101011101010101)$ and $x_2 = 1$, then $H(x_1) = (101)$ and $H(x_2) = (001)$. It is easy to see that this hash function is deterministic and uniform.

Of particular interest are so-called **cryptographic** hash functions, which are hash functions that are also **one-way functions**, which essentially means that, given a string y from $\{0,1\}^k$ it is practically infeasible to find a string $x \in \{0,1\}^*$ such that H(x) = y holds. This property is usually called **preimage-resistance**.

are these going to be relevant later?

TODO: theorem: every factor of order defines a subgroup...

In addition, it should be infeasible to find two strings $x_1, x_2 \in \{0, 1\}^*$, such that $H(x_1) = H(x_2)$, which is called **collision resistance**. It is important to note, though, that collisions always exist, since a function $H: \{0, 1\}^* \to \{0, 1\}^k$ inevitably maps infinitely many values onto the same hash. In fact, for any hash function with digests of length k, finding a preimage to a given digest can always be done using a brute force search in 2^k evaluation steps. It should just be practically impossible to compute those values, and statistically very unlikely to generate two of them by chance.

A third property of a cryptographic hash function is that small changes in the input string, like changing a single bit, should generate hash values that look completely different from each other.

Because cryptographically secure hash functions map tiny changes in input values onto large changes in the output, implementation errors that change the outcome are usually easy to spot by comparing them to expected output values. The definitions of cryptographically secure hash functions are therefore usually accompanied by some test vectors of common inputs and expected digests. Since the empty string " is the only string of length 0, a common test vector is the expected digest of the empty string.

Example 44 (k-truncation hash). Consider the k-truncation hash from example 43. Since the empty string has length 0, it follows that the digest of the empty string is string of length k that only contains 0's:

$$H_k('') = (000...000)$$
 (4.11)

It is pretty obvious from the definition of H_k that this simple hash function is not a cryptographic hash function. In particular, every digest is its own preimage, since $H_k(y) = y$ for every string of size exactly k. Finding preimages is therefore easy, so the property of preimage resistance does not hold.

In addition, it is easy to construct collisions as all strings of size > k that share the same k-bits on the right are mapped to the same hash value, so this function is not collision resistant, either

Finally, this hash function is not very chaotic, as changing bits that are not part of the k right-most bits don't change the digest at all.

Computing cryptographically secure hash functions in pen-and-paper style is possible but tedious. Fortunately, Sage can import the **hashlib** library, which is intended to provide a reliable and stable base for writing Python programs that require cryptographic functions. The following examples explain how to use hashlib in Sage.

Example 45. An example of a hash function that is generally believed to be a cryptographically secure hash function is the so-called **SHA256** hash, which, in our notation, is a function that maps binary strings of arbitrary length onto binary strings of length 256:

$$SHA256: \{0,1\}^* \to \{0,1\}^{256}$$
 (4.12)

To evaluate a proper implementation of the *SHA*256 hash function, the digest of the empty string is supposed to be

$$SHA256('') = e3b0c44298fc1c149afbf4c8996fb92427ae41e4649b934ca495991b7852b855$$
 (4.13)

For better human readability, it is common practice to represent the digest of a string not in its binary form, but in a hexadecimal representation. We can use Sage to compute *SHA*256 and freely transit between binary, hexadecimal and decimal representations. To do so, we import hashlib's implementation of SHA256:

Is there a term for this property?

```
sage: import hashlib
                                                         144
1988
   sage: test = 'e3b0c44298fc1c149afbf4c8996fb92427ae41e4649b934
                                                         145
1989
     ca495991b7852b855'
1990
   sage: hasher = hashlib.sha256(b'')
                                                         146
1991
   sage: str = hasher.hexdigest()
                                                         147
1992
   sage: type(str)
                                                         148
1993
   <class 'str'>
                                                         149
1994
   sage: d = ZZ('0x' + str) \# conversion to integer type
                                                         150
1995
   sage: d.str(16) == str
                                                         151
1996
   True
                                                         152
1997
   sage: d.str(16) == test
                                                         153
1998
                                                         154
1999
   sage: d.str(16)
                                                         155
2000
   e3b0c44298fc1c149afbf4c8996fb92427ae41e4649b934ca495991b7852b8
                                                         156
2001
2002
   sage: d.str(2)
                                                         157
2003
   158
2004
     2005
     2006
     2007
     01011100001010101
2008
   sage: d.str(10)
                                                         159
2009
   10298733624955409702953521232258132278979990064819803499337939
                                                         160
2010
     7001115665086549
2011
```

Hashing to cyclic groups As we have seen in the previous paragraph, general hash functions map binary strings of arbitrary length onto binary strings of length k. However, it is desirable in various cryptographic primitives to not simply hash to binary strings of fixed length but to hash into algebraic structures like groups, while keeping (some of) the properties like preimage resistance or collision resistance.

Hash functions like this can be defined for various algebraic structures, but, in a sense, the most fundamental ones are hash functions that map into groups, because they can be easily extended to map into other structures like rings or fields.

To give a more precise definition, let \mathbb{G} be a group and $\{0,1\}^*$ the set of all finite, binary strings, then a **hash-to-group** function is a deterministic map

$$H: \{0,1\}^* \to \mathbb{G} \tag{4.14}$$

Common properties of hash functions, like uniformity, are desirable but not always realized in real-world instantiations of hash-to-group functions, so we skip those requirements for now and keep the definition very general.

As the following example shows, hashing to finite cyclic groups can be trivially achieved for the price of some undesirable properties of the hash function:

Example 46 (Naive cyclic group hash). Let \mathbb{G} be a finite cyclic group. If the task is to implement a hash-to-group function, one immediate approach can be based on the observation that binary strings of size k can be interpreted as integers $z \in \mathbb{Z}$ in the range $0 \le z < 2^k$.

To be more precise, choose an ordinary hash function $H: \{0,1\}^* \to \{0,1\}^k$ for some parameter k and a generator g of \mathbb{G} . Then the expression below is a positive integer (where $H(s)_i$

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means the bit at the *j*-th position of H(s):

$$z_{H(s)} = H(s)_0 \cdot 2^0 + H(s)_1 \cdot 2^1 + \dots + H(s)_k \cdot 2^k$$
(4.15)

A hash-to-group function for the group \mathbb{G} can then be defined as a concatenation of the exponential map $g^{(\cdot)}$ of g with the interpretation of H(s) as an integer:

$$H_g: \{0,1\}^* \to \mathbb{G}: s \mapsto g^{z_{H(s)}}$$
 (4.16)

Constructing a hash-to-group function like this is easy for cyclic groups, and it might be good enough in certain applications. It is, however, almost never adequate in cryptographic applications, as discrete log relations might be constructible between two given hash values $H_g(s)$ and $H_g(t)$.

a few examples?

To see that, assume that \mathbb{G} is of order r and that $z_{H(s)}$ has a multiplicative inverse in modular r arithmetics. In that case, we can compute $x = z_{H(t)} \cdot z_{H(s)}^{-1}$ in \mathbb{Z}_r and find a discrete log relation between the group hash values, that is, find some x with $H_g(t) = (H_g(s))^x$:

$$H_g(t) = (H_g(s))^x \Leftrightarrow g^{z_{H(t)}} = g^{z_{H(t)}} \Leftrightarrow g^{z_{H(t)}}$$

Therefore applications where discrete log relations between hash values are undesirable need different approaches. Many of these approaches start with a way to hash into the set \mathbb{Z}_r of modular r arithmetics.

Hashing to modular arithmetics One of the most widely used applications of hash-intogroup functions are hash functions that map into the set \mathbb{Z}_r of modular r arithmetics for some modulus r. Different approaches to construct such a function are known, but probably the most widely used ones are based on the insight that the images of arbitrary hash functions can be interpreted as binary representations of integers, as explained in example 46.

check reference

It follows from this interpretation that one simple method of hashing into \mathbb{Z}_r is constructed by observing that if r is a modulus with a bit-length of k = r.nbits(), then every binary string $(b_0, b_1, \ldots, b_{k-2})$ of length k-1 defines an integer z in the rage $0 \le z < 2^{k-1} \le r$, by defining

$$z = b_0 \cdot 2^0 + b_1 \cdot 2^1 + \dots + b_{k-2} \cdot 2^{k-2}$$
(4.17)

Now, since z < r, we know that z is guaranteed to be in the set $\{0, 1, \dots, r-1\}$, and hence it can be interpreted as an element of \mathbb{Z}_r . From this it follows that if $H: \{0, 1\}^* \to \{0, 1\}^{k-1}$ is a hash function, then a hash-to-group function can be constructed as follows (where $H(s)_j$ means the j-th bit of the image binary string H(s) of the original binary hash function):

$$H_{r,nbits()-1}: \{0,1\}^* \to \mathbb{Z}_r: s \mapsto H(s)_0 \cdot 2^0 + H(s)_1 \cdot 2^1 + \dots + H(s)_{k-2} \cdot 2^{k-2}$$
 (4.18)

A drawback of this hash function is that the distribution of the hash values in \mathbb{Z}_r is not necessarily uniform. In fact, if $r-2^{k-1} \neq 0$, then by design $H_{r,nbits()-1}$ will never hash onto values $z \geq 2^{k-1}$. Good moduli r are therefore as close to 2^{k-1} as possible, while less good moduli are closer to 2^k . In the worst case, when $r = 2^k - 1$, it misses $2^{k-1} - 1$, that is, almost half of all elements, from \mathbb{Z}_r .

An advantage of this approach is that properties like preimage resistance or collision resistance of the original hash function $H(\cdot)$ are preserved.

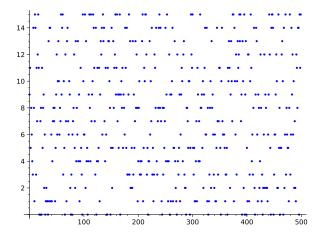
TODO: DOUBLE CHECK THIS REA-SONING. *Example* 47. To give an implementation of the $H_{r,nbits()-1}$ hash function, we use a 5-bit truncation of the SHA256 hash from example 45 and define a hash into \mathbb{Z}_{16} as follows:

$$H_{16.nbits()-1}: \{0,1\}^* \to \mathbb{Z}_{16}: s \mapsto SHA256(s)_0 \cdot 2^0 + SHAH256(s)_1 \cdot 2^1 + \ldots + SHA256(s)_4 \cdot 2^4 + \ldots + SHA256(s)_6 \cdot 2^6 + SHAH256(s)_6 \cdot 2^6 + S$$

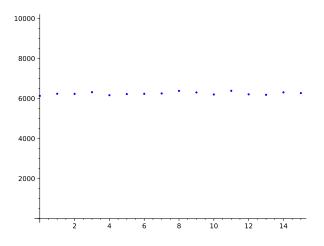
Since k = 16.nbits() = 5 and $16 - 2^{k-1} = 0$, this hash maps uniformly onto \mathbb{Z}_{16} . We can invoke Sage to implement it:

```
sage: import hashlib
                                                                                   161
2064
    sage: def Hash5(x):
                                                                                   162
2065
                 hasher = hashlib.sha256(x)
                                                                                   163
2066
                 digest = hasher.hexdigest()
                                                                                   164
2067
                 d = ZZ(digest, base=16)
                                                                                   165
2068
                 d = d.str(2)[-4:]
                                                                                   166
2069
                 return ZZ(d,base=2)
                                                                                   167
2070
    sage: Hash5(b'')
                                                                                   168
2071
    5
                                                                                   169
2072
```

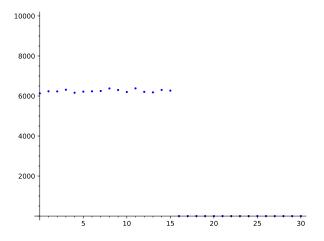
We can then use Sage to apply this function to a large set of input values in order to plot a visualization of the distribution over the set $\{0, ..., 15\}$. Executing over 500 input values gives the following plot:



To get an intuition of uniformity, we can count the number of times the hash function $H_{16.nbits()-1}$ maps onto each number in the set $\{0,1,\ldots,15\}$ in a loop of 100000 hashes, and compare that to the ideal uniform distribution, which would map exactly 6250 samples to each element. This gives the following result:



The uniformity of distribution problem becomes apparent if we want to construct a similar hash function for \mathbb{Z}_r for any r in the range $17 \le r \le 31$. In this case, the definition of the hash function is exactly the same as for \mathbb{Z}_{16} , and hence, the images will not exceed the value 16. So, for example, even in the case of hashing to \mathbb{Z}_{31} , the hash function never maps to any value larger than 16, leaving almost half of all numbers out of the image range.



The second widely used method of hashing into \mathbb{Z}_r is constructed by observing the following: If r is a modulus with a bit-length of $r.bits() = k_1$ and $H : \{0,1\}^* \to \{0,1\}^{k_2}$ is a hash function that produces digests of size k_2 , with $k_2 \ge k_1$, then a hash-to-group function can be constructed by interpreting the image of H as a binary representation of an integer and then taking the modulus by r to map into \mathbb{Z}_r . This is formalized in the equation below, where $H(s)_j$ means the j'th bit of the image binary string H(s) of the original binary hash function.

Mirco: We can do better than this

$$H'_{mod_r}: \{0,1\}^* \to \mathbb{Z}_r: s \mapsto \left(H(s)_0 \cdot 2^0 + H(s)_1 \cdot 2^1 + \dots + H(s)_{k_2} \cdot 2^{k_2}\right) \bmod n$$
 (4.19)

A drawback of this hash function is that computing the modulus requires some computational effort. In addition, the distribution of the hash values in \mathbb{Z}_r might not be even, depending on the difference $2^{k_2+1}-r$. An advantage of it is that potential properties of the original hash function $H(\cdot)$ (like preimage resistance or collision resistance) are preserved, and the distribution can be made almost uniform, with only negligible bias depending on what modulus r and images size k_2 are chosen.

Example 48. To give an implementation of the H_{mod_r} hash function, we use k_2 -bit truncation of the SHA256 hash from example 45, and define a hash into \mathbb{Z}_{23} as follows:

$$H_{mod_{23},k_2}: \{0,1\}^* \to \mathbb{Z}_{23}:$$

$$s \mapsto \left(SHA256(s)_0 \cdot 2^0 + SHAH256(s)_1 \cdot 2^1 + \ldots + SHA256(s)_{k_2} \cdot 2^{k_2}\right) \mod 23$$

We want to use various instantiations of k_2 to visualize the impact of truncation length on the distribution of the hashes in \mathbb{Z}_{23} . We can invoke Sage to implement it as follows:

```
sage: import hashlib
                                                                               170
2103
    sage: Z23 = Integers(23)
                                                                               171
2104
    sage:
           def Hash_mod23(x, k2):
                                                                               172
2105
                hasher = hashlib.sha256(x.encode('utf-8'))
                                                                               173
2106
                digest = hasher.hexdigest()
                                                                               174
2107
                d = ZZ(digest, base=16)
2108
                                                                               175
```

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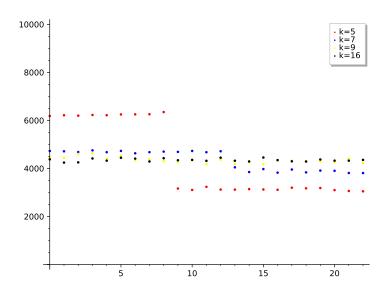
2126

```
2109 ....: d = d.str(2)[-k2:] 176

2110 ....: d = ZZ(d, base=2) 177

2111 ....: return Z23(d) 178
```

We can then use Sage to apply this function to a large set of input values in order to plot visualizations of the distribution over the set $\{0, \dots, 22\}$ for various values of k_2 , by counting the number of times it maps onto each number in a loop of 100000 hashes. We get the following plot:



A third method that can sometimes be found in implementations is the so-called "**try-and-increment**" **method**. To understand this method, we define an integer $z \in \mathbb{Z}$ from any hash value H(s) as we did in the previous methods, that is, we define $z = H(s)_0 \cdot 2^0 + H(s)_1 \cdot 2^1 + \dots + H(s)_{k-1} \cdot 2^k$.

Hashing into \mathbb{Z}_r is then achievable by first computing z, and then trying to see if $z \in \mathbb{Z}_r$. If it is, then the hash is done; if not, the string s is modified in a deterministic way and the process is repeated until a suitable number z is found. A suitable, deterministic modification could be to concatenate the original string by some bit counter. A "try-and-increment" algorithm would then work like in algorithm 5.

check reference

```
Algorithm 5 Hash-to-\mathbb{Z}_n
```

```
Require: r \in \mathbb{Z} with r.nbits() = k and s \in \{0,1\}^*

procedure TRY-AND-INCREMENT(r,k,s)

c \leftarrow 0

repeat

s' \leftarrow s||c\_bits()

z \leftarrow H(s')_0 \cdot 2^0 + H(s')_1 \cdot 2^1 + \ldots + H(s')_k \cdot 2^k

c \leftarrow c + 1

until z < r

return x

end procedure

Ensure: z \in \mathbb{Z}_r
```

Depending on the parameters, this method can be very efficient. In fact, if k is sufficiently large and r is close to 2^{k+1} , the probability for z < r is very high and the repeat loop will almost

always be executed a single time only. A drawback is, however, that the probability of having to execute the loop multiple times is not zero.

Once some hash function into modular arithmetics is found, it can often be combined with additional techniques to hash into more general finite cyclic groups. The following paragraphs describe a few of those methods widely adopted in SNARK development.

Pedersen Hashes The so-called **Pedersen hash function** [Pedersen, 1992] provides a way to map binary inputs of fixed size *k* onto elements of finite cyclic groups that avoids discrete log relations between the images as they occur in the naive approach XXX. Combining it with a classical hash function provides a hash function that maps strings of arbitrary length onto group elements.

add reference

To be more precise, let j be an integer, \mathbb{G} a finite cyclic group of order r and $\{g_1, \ldots, g_j\} \subset \mathbb{G}$ a uniform randomly generated set of generators of \mathbb{G} . Then **Pedersen's hash function** is defined as follows:

$$H_{Ped}: (\mathbb{Z}_r)^j \to \mathbb{G}: (x_1, \dots, x_j) \mapsto \prod_{i=1}^j g_i^{x_j}$$
 (4.20)

It can be shown that Pedersen's hash function is collision-resistant under the assumption that \mathbb{G} is a DL-A group. However, it is important to note that Pedersen hashes cannot be assumed to be pseudorandom and should therefore not be used where a hash function serves as an approximation of a random oracle, will these be explained in the initial chapters?

pseudorandor

oracle

From an implementation perspective, it is important to derive the set of generators $\{g_1, \ldots, g_j\}$ in such a way that they are as uniform and random as possible. In particular, any known discrete log relation between two generators, that is, any known $x \in \mathbb{Z}_r$ with $g_h = (g_i)^x$ must be avoided.

To see how Pedersen hashes can be used to define an actual hash-to-group function according to our definition, we can use any of the hash-to- \mathbb{Z}_r functions as we have derived them in equation 4.18.

check reference

MimC Hashes [Albrecht et al., 2016]

add text

Pseudorandom Functions in DDH-A groups As noted in above, Pederson's hash function does not have the properties a random function and should therefore not be instantiated as such. To see an example of a random oracle function in groups where the decisional Diffie-Hellman construction is assumed to hold true, let \mathbb{G} be a DDH-A group of order r with generator g and $\{a_0, a_1, \ldots, a_k\} \subset \mathbb{Z}_r^*$ a uniform randomly generated set of numbers invertible in modular r arithmetics. Then a pseudo-random function is given by the as follows:

$$F_{rand}: \{0,1\}^{k+1} \to \mathbb{G}: (b_0,\dots,b_k) \mapsto g^{b_0 \cdot \prod_{i=1}^k a_i^{b_i}}$$
 (4.21)

Of course, if $H: \{0,1\}^* \to \{0,1\}^{k+1}$ is a random oracle, then the concatenation of F_{rand} and H_{2159} also defines a random oracle

$$H_{rand,\mathbb{G}}: \{0,1\}^* \to \mathbb{G}: s \mapsto F_{rand}(H(s))$$
 (4.22)

4.2 Commutative Rings

Thinking back to operations on integers, we know that there are two of these: addition and multiplication. As we have seen, addition defines a group structure on the set of integers.

However, multiplication does not define a group structure, given that integers generally don't have multiplicative inverses.

Configurations like these constitute so-called **commutative rings with unit**. To be more precise, a commutative ring with unit $(R, +, \cdot, 1)$ is a set R provided with two maps $+: R \cdot R \to R$ and $\cdot: R \cdot R \to R$, called **addition** and **multiplication**, such that the following conditions hold:

2168 Definition 4.2.0.1. Commutative ring with unit

- (R,+) is a commutative group, where the neutral element is denoted with 0. Commutativity of multiplication: $r_1 \cdot r_2 = r_2 \cdot r_1$ for all $r_1, r_2 \in R$.
- Existence of a unit: There is an element $1 \in R$, such that $1 \cdot g$ holds for all $g \in R$,
- Associativity: For every $g_1, g_2, g_3 \in \mathbb{G}$ the equation $g_1 \cdot (g_2 \cdot g_3) = (g_1 \cdot g_2) \cdot g_3$ holds.
- **Distributivity**: For all $g_1, g_2, g_3 \in R$ the distributive laws $g_1 \cdot (g_2 + g_3) = g_1 \cdot g_2 + g_1 \cdot g_3$ holds.

Example 49 (The ring of integers). The set \mathbb{Z} of integers with the usual addition and multiplication is the archetypical example of a commutative ring with unit 1.

Example 50 (Underlying commutative group of a ring). Every commutative ring with unit $(R, +, \cdot, 1)$ gives rise to a group, if we disregard multiplication.

The following example is somewhat unusual, but we encourage you to think through it because it helps to detach the mind from familiar styles of computation and concentrate on the abstract algebraic explanation.

Example 51. Let $S := \{ \bullet, \star, \odot, \otimes \}$ be a set that contains four elements, and let addition and multiplication on S be defined as follows:

Then (S, \cup, \circ) is a ring with unit \star and zero \bullet . It therefore makes sense to ask for solutions to equations like this one: Find $x \in S$ such that

$$\otimes \circ (x \cup \odot) = \star$$

To see how such a "moonmath equation" can be solved, we have to keep in mind that rings behaves mostly like normal numbers when it comes to bracketing and computation rules. The only differences are the symbols, and the actual way to add and multiply them. With this in

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mind, we solve the equation for x in the "usual way":

So, even though this equation looked really alien at first glance, we could solve it basically exactly the way we solve "normal" equations containing numbers.

Note, however, that whenever a multiplicative inverse would be needed to solve an equation in the usual way in a ring, things can be very different than most of us are used to. For example, the following equation cannot be solved for x in the usual way, since there is no multiplicative inverse for \odot in our ring.

$$\odot \circ x = \otimes \tag{4.23}$$

We can confirm this by looking at the multiplication table to see that no such x exits. As another example, the following equation does not have a single solution but two: $x \in \{\star, \otimes\}$.

$$\bigcirc \circ x = \bigcirc$$
 (4.24)

Having no solution or two solutions is certainly not something to expect from types like \mathbb{Q} (rational numbers). can we use another set as an example? we hardly talked about Q so far

Example 52. Considering polynomials again, we note from their definition that what we have called the type R of the coefficients must in fact be a commutative ring with a unit, since we need addition, multiplication, commutativity and the existence of a unit for R[x] to have the properties we expect.

Considering R to be a ring with addition and multiplication of polynomials as defined in 4.2.0.1 actually makes R[x] into a commutative ring with a unit, too, where the polynomial 1 is the multiplicative unit.

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Example 53. Let n be a modulus and $(\mathbb{Z}_n, +, \cdot)$ the set of all remainder classes of integers modulo n, with the projection of integer addition and multiplication as defined in 4.2.0.1. It can be shown that $(\mathbb{Z}_n, +, \cdot)$ is a commutative ring with unit 1.

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Considering the exponential map from page 43 <u>again, let \mathbb{G} be a finite cyclic group of order</u> n with generator $g \in \mathbb{G}$. Then the ring structure of $(\mathbb{Z}_n, +, \cdot)$ is mapped onto the group structure of \mathbb{G} in the following way:

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$$g^{x+y} = g^x \cdot g^y$$
 for all $x, y \in \mathbb{Z}_n$
 $g^{x \cdot y} = (g^x)^y$ for all $x, y \in \mathbb{Z}_n$

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This is of particular interest in cryptography and SNARKs, as it allows for the evaluation of polynomials with coefficients in \mathbb{Z}_n to be evaluated "in the exponent". To be more precise, let $p \in \mathbb{Z}_n[x]$ be a polynomial with $p(x) = a_m \cdot x^m + a_{m-1}x^{m-1} + \ldots + a_1x + a_0$. Then the previously defined exponential laws (example 37) imply the following:

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$$g^{p(x)} = g^{a_m \cdot x^m + a_{m-1} x^{m-1} + \dots + a_1 x + a_0}$$

= $\left(g^{x^m} \right)^{a_m} \cdot \left(g^{x^{m-1}} \right)^{a_{m-1}} \cdot \dots \cdot (g^x)^{a_1} \cdot g^{a_0}$

Hence, to evaluate p at some point s in the exponent, we can insert s into the right-hand side of the last equation and evaluate the product.

As we will see, this is a key insight to understanding many SNARK protocols like e.g. Groth16 [Groth, 2016] or XXX.

Example 54. To give an example of the evaluation of a polynomial in the exponent of a finite cyclic group, consider the exponential map from example 37:

$$3^{(\cdot)}: \mathbb{Z}_4 \to \mathbb{F}_5^* x \mapsto 3^x$$

Choosing the polynomial $p(x) = 2x^2 + 3x + 1$ from $\mathbb{Z}_4[x]$, we can evaluate the polynomial at say x = 2 in the exponent of 3 in two different ways. On the one hand, we can evaluate p at 2 and then write the result into the exponent as follows:

$$3^{p(2)} = 3^{2 \cdot 2^2 + 3 \cdot 2 + 1}$$

$$= 3^{2 \cdot 0 + 2 + 1}$$

$$= 3^3$$

$$= 2$$

On the other hand, we can use the right-hand side of the equation to evaluate p at 2 in the exponent of 3 as follows:

$$3^{p(2)} = \left(3^{2^2}\right)^2 \cdot \left(3^2\right)^3 \cdot 3^1$$

$$= \left(3^0\right)^2 \cdot 3^3 \cdot 3$$

$$= 1^2 \cdot 2 \cdot 3$$

$$= 2 \cdot 3$$

$$= 2$$

Hashing to Commutative Rings As we have seen in XXX, various constructions for hashing-to-groups are known and used in applications. As commutative rings are Abelian groups, when we disregard the multiplicative structure, hash-to-group constructions can be applied for hashing into commutative rings, too. This is possible in general, as the codomain of a general hash function $\{0,1\}^*$ is just the set of binary strings of arbitrary but finite length, which has no algebraic structure that the hash function must respect.

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Abelian groups

codomain

4.3 Fields

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We started this chapter with the definition of a group (section 4.1), which we then expanded into the definition of a commutative ring with a unit (section 4.2). Such rings generalize the behavior

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of integers. In this section, we will look at those special cases of commutative rings where every element other than the neutral element of addition has a multiplicative inverse. Those structures behave very much like the rational numbers Q. Rational numbers are, in a sense, an extension of the ring of integers, that is, they are constructed by including newly defined multiplicative inverses (fractions) to the integers.

Now, considering the definition of a ring (4.2.0.1) again, we define a **field** $(\mathbb{F},+,\cdot)$ to be a set \mathbb{F} , together with two maps $+: \mathbb{F} \cdot \mathbb{F} \to \mathbb{F}$ and $\cdot: \mathbb{F} \cdot \mathbb{F} \to \mathbb{F}$, called *addition* and *multiplication*, such that the following conditions hold:

Definition 4.3.0.1. Field

- $(\mathbb{F},+)$ is a commutative group, where the neutral element is denoted by 0.
- $(\mathbb{F}\setminus\{0\},\cdot)$ is a commutative group, where the neutral element is denoted by 1.
- (Distributivity) For all $g_1, g_2, g_3 \in \mathbb{F}$ the distributive law $g_1 \cdot (g_2 + g_3) = g_1 \cdot g_2 + g_1 \cdot g_3$ holds.

If a field is given and the definition of its addition and multiplication is not ambiguous, we will often simply write \mathbb{F} instead of $(\mathbb{F}, +, \cdot)$ to denote the field. Moreover, we use \mathbb{F}^* to describe the multiplicative group of the field, that is, the set of elements with multiplication as the group law, excluding the neutral element of addition.

The **characteristic** of a field \mathbb{F} , represented as $char(\mathbb{F})$, is the smallest natural number $n \ge 1$ for which the *n*-fold sum of 1 equals zero, i.e. for which $\sum_{i=1}^{n} 1 = 0$. If such an n > 0 exists, the field is also said to have a **finite characteristic**. If, on the other hand, every finite sum of 1 is such that it is not equal to zero, then the field is defined to have characteristic 0. S: Tried to disambiguate the scope of negation between 1. "It is true of every finite sum of 1 that it is not equal to 0" and 2. "It is not true of every finite sum of 1 that it is equal to 0" From the example below, it looks like 1. is the intended meaning here, correct?

Check change of wording

Example 55 (Field of rational numbers). Probably the best known example of a field is the set of rational numbers Q together with the usual definition of addition, subtraction, multiplication and division. Since there is no natural number $n \in \mathbb{N}$, such that $\sum_{i=0}^{n} 1 = 0$ in the set of rational numbers, the characteristic $char(\mathbb{Q})$ of the field \mathbb{Q} is zero. In Sage, rational numbers are called as follows:

```
sage: QQ
                                                                                  179
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    Rational Field
                                                                                  180
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    sage: QQ(1/5) # Get an element from the field of rational
                                                                                  181
        numbers
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    1/5
                                                                                  182
2252
    sage: QQ(1/5) / QQ(3) \# Division
                                                                                  183
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    1/15
                                                                                  184
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```

Example 56 (Field with two elements). It can be shown that, in any field, the neutral element 0 of addition must be different from the neutral element 1 of multiplication, that is, $0 \neq 1$ always holds in a field. From this, it follows that the smallest field must contain at least two elements. As the following addition and multiplication tables show, there is indeed a field with two elements, which is usually called \mathbb{F}_2 :

Let $\mathbb{F}_2 := \{0,1\}$ be a set that contains two elements and let addition and multiplication on 2260 \mathbb{F}_2 be defined as follows:

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Since 1+1=0 in the field \mathbb{F}_2 , we know that the characteristic of \mathbb{F}_2 is there, that is, we have 2263 $char(\mathbb{F}_2) = 0.$ 2264

For reasons we will understand better in XXX, Sage defines this field as a so-called Galois field with 2 elements. You can call it in Sage as follows:

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```
sage: F2 = GF(2)
                                                                                     185
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     sage: F2(1) # Get an element from GF(2)
                                                                                     186
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                                                                                     187
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     sage: F2(1) + F2(1) \# Addition
                                                                                     188
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                                                                                     189
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     sage: F2(1) / F2(1) # Division
                                                                                     190
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                                                                                     191
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```

Example 57. Both the real numbers \mathbb{R} as well as the complex numbers \mathbb{C} are well known examples of fields. 2275

Expand on this?

Exercise 30. Consider our remainder class ring $(\mathbb{F}_5,+,\cdot)$ and show that it is a field. What is the 2276 characteristic of \mathbb{F}_5 ? 2277

Prime fields As we have seen in the various examples of the previous sections, modular arithmetics behaves similarly to the ordinary arithmetics of integers in many ways. This is due to the fact that remainder class sets \mathbb{Z}_n are commutative rings with units.

However, we have also seen in 36 that, whenever the modulus is a prime number, every check remainder class other than the zero class has a modular multiplicative inverse. This is an important observation, since it immediately implies that, in case of a prime number, the remainder class set \mathbb{Z}_n is not just a ring but actually a **field**. Moreover, since $\sum_{i=0}^n 1 = 0$ in \mathbb{Z}_n , we know that those fields have the finite characteristic n.

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To distinguish this important case from arbitrary remainder class rings, we write $(\mathbb{F}_p,+,\cdot)$ for the field of all remainder classes for a prime number modulus $p \in \mathbb{P}$ and call it the **prime field** of characteristic p.

Prime fields are the foundation for many of the contemporary algebra-based cryptographic systems, as they have many desirable properties. One of them is that, since these sets are finite and a prime field of characteristic, p can be represented on a computer in roughly $log_2(p)$ amount of space without precision problems that are unavoidable for computer representations of infinite sets such as rational numbers or integers.

Since prime fields are special cases of remainder class rings, all computations remain the same. Addition and multiplication can be computed by first doing normal integer addition and multiplication, and then taking the remainder modulus p. Subtraction and division can be computed by adding or multiplying with the additive or the multiplicative inverse, respectively. The additive inverse -x of a field element $x \in \mathbb{F}_p$ is given by p-x, and the multiplicative inverse of $x \neq 0$ is given by x^{p-2} , or can be computed using the Extended Euclidean Algorithm.

Note that these computations might not be the fastest to implement on a computer. They are, however, useful in this book, as they are easy to compute for small prime numbers.

Example 58. The smallest field is the field \mathbb{F}_2 of characteristic 2 as we have seen in example 56. It is the prime field of the prime number 2.

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Example 59. To summarize the basic aspects of computation in prime fields, let us consider the prime field \mathbb{F}_5 and simplify the following expression:

$$\left(\frac{2}{3}-2\right)\cdot 2$$

A first thing to note is that since \mathbb{F}_5 is a field, all rules are identical to the rules we learned in school when we where dealing with rational, real or complex numbers. This means we can use e.g. bracketing (distributivity) or addition as usual:

$$\left(\frac{2}{3}-2\right) \cdot 2 = \frac{2}{3} \cdot 2 - 2 \cdot 2$$

$$= \frac{2 \cdot 2}{3} - 2 \cdot 2$$

$$= \frac{4}{3} - 4$$

$$= 4 \cdot 2 - 4$$

$$= 4 \cdot 2 + 1$$

$$= 3 + 1$$

$$= 4$$
distributive law
$$4 \mod 5 = 4$$
additive inverse of 3 is $3^{5-2} \mod 5 = 2$
additive inverse of 4 is $5 - 4 = 1$

$$8 \mod 5 = 3$$

$$4 \mod 5 = 4$$

In this computation, we computed the multiplicative inverse of 3 using the identity $x^{-1} = x^{p-2}$ in a prime field. This is impractical for large prime numbers. Recall that another way of computing the multiplicative inverse is the Extended Euclidean Algorithm (see 3.10 on page 18). To refresh our memory, the task is to compute $x^{-1} \cdot 3 + t \cdot 5 = 1$, but t is actually irrelevant. We get

So the multiplicative inverse of 3 in \mathbb{Z}_5 is 2 and indeed if compute $3 \cdot 2$ we get 1 in \mathbb{F}_5 .

Square Roots In this part, we deal with square numbers, also called **quadratic residues** and **square roots** in prime fields. This is of particular importance in our studies on elliptic curves, as only square numbers can actually be points on an elliptic curve.

To make the intuition of quadratic residues and roots precise, let $p \in \mathbb{P}$ be a prime number and \mathbb{F}_p its associated prime field. Then a number $x \in \mathbb{F}_p$ is called a **square root** of another number $y \in \mathbb{F}_p$, if x is a solution to the following equation:

$$x^2 = y \tag{4.25}$$

In this case, y is called a **quadratic residue**. On the other hand, if y is given and the quadratic equation has no solution x, we call y a **quadratic non-residue**. For any $y \in \mathbb{F}_p$, we denote the set of all square roots of y in the prime field \mathbb{F}_n as follows:

$$\sqrt{y} := \{ x \in \mathbb{F}_p \mid x^2 = y \}$$
 (4.26)

unify \mathbb{Z}_5 and \mathbb{F}_5 across the board?

S: are we introducing elliptic curves in section 1 or 2?

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If y is a quadratic non-residue, then $\sqrt{y} = \emptyset$ (an empty set), and if y = 0, then $\sqrt{y} = \{0\}$.

Informally speaking, quadratic residues are numbers such that we can take the square root of them, while quadratic non-residues are numbers that don't have square roots. The situation therefore parallels the familiar case of integers, where some integers like 4 or 9 have square roots and others like 2 or 3 don't (as integers).

It can be shown that, in any prime field, every non zero element has either no square root or two of them. We adopt the convention to call the smaller one (when interpreted as an integer) as the **positive** square root and the larger one as the **negative**. This makes sense, as the larger one can always be computed as the modulus minus the smaller one, which is the definition of the negative in prime fields.

Example 60 (Quadratic (Non)-Residues and roots in \mathbb{F}_5). Let us consider our example prime field \mathbb{F}_5 again. All square numbers can be found on the main diagonal of the multiplication table in example 13 on page 28. As you can see, in \mathbb{F}_5 only the numbers 0, 1 and 4 have square roots and we get $\sqrt{0} = \{0\}$, $\sqrt{1} = \{1,4\}$, $\sqrt{2} = \emptyset$, $\sqrt{3} = \emptyset$ and $\sqrt{4} = \{2,3\}$. The numbers 0, 1 and 4 are therefore quadratic residues, while the numbers 2 and 3 are quadratic non-residues.

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In order to describe whether an element of a prime field is a square number or not, the so-called **Legendre symbol** can sometimes be found in the literature, which is why we will summerize it here:

Let $p \in \mathbb{P}$ be a prime number and $y \in \mathbb{F}_p$ an element from the associated prime field. Then the *Legendre symbol* of y is defined as follows:

$$\left(\frac{y}{p}\right) := \begin{cases} 1 & \text{if } y \text{ has square roots} \\ -1 & \text{if } y \text{ has no square roots} \\ 0 & \text{if } y = 0 \end{cases}$$
 (4.27)

Example 61. Looking at the quadratic residues and non residues in \mathbb{F}_5 from example 13 again, we can deduce the following Legendre symbols, from example XXX.

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$$\left(\frac{0}{5}\right) = 0$$
, $\left(\frac{1}{5}\right) = 1$, $\left(\frac{2}{5}\right) = -1$, $\left(\frac{3}{5}\right) = -1$, $\left(\frac{4}{5}\right) = 1$.

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The Legendre symbol provides a criterion to decide whether or not an element from a prime field has a quadratic root or not. This, however, is not just of theoretical use: The so-called **Euler criterion** provides a compact way to actually compute the Legendre symbol. To see that, let $p \in \mathbb{P}_{\geq 3}$ be an odd prime number and $y \in \mathbb{F}_p$. Then the Legendre symbol can be computed as follows:

$$\left(\frac{y}{p}\right) = y^{\frac{p-1}{2}} \,. \tag{4.28}$$

Example 62. Looking at the quadratic residues and non residues in \mathbb{F}_5 from example 13 again,

check reference we can compute the following Legendre symbols using the Euler criterion:

$$\left(\frac{0}{5}\right) = 0^{\frac{5-1}{2}} = 0^2 = 0$$

$$\left(\frac{1}{5}\right) = 1^{\frac{5-1}{2}} = 1^2 = 1$$

$$\left(\frac{2}{5}\right) = 2^{\frac{5-1}{2}} = 2^2 = 4 = -1$$

$$\left(\frac{3}{5}\right) = 3^{\frac{5-1}{2}} = 3^2 = 4 = -1$$

$$\left(\frac{4}{5}\right) = 4^{\frac{5-1}{2}} = 4^2 = 1$$

Exercise 31. Consider the prime field \mathbb{F}_{13} . Find the set of all pairs $(x,y) \in \mathbb{F}_{13} \times \mathbb{F}_{13}$ that satisfy the equation

$$x^2 + y^2 = 1 + 7 \cdot x^2 \cdot y^2$$

Exponentiation TO APPEAR...

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Hashing into prime fields An important problem in SNARK development is the ability to hash to (various subsets) of elliptic curves. As we will see in XXX, those curves are often defined over prime fields, and hashing to a curve might start with hashing to the prime field. It is therefore important to understand how to hash into prime fields.

On pages 51–55, we looked at a few methods of hashing into the residue class rings \mathbb{Z}_n for arbitrary n > 1. As prime fields are just special instances of those rings, all methods for hashing into \mathbb{Z}_n functions can be used for hashing into prime fields, too.

Extension Fields Prime fields, defined in the previous section, are the basic building blocks for cryptography in general and SNARKs in particular.

However, as we will see in XXX so-called **pairing-based** SNARK systems are crucially dependent on group pairings XXX defined over the group of rational points of elliptic curves. For those pairings to be non-trivial, the elliptic curve must not only be defined over a prime field, but over a so-called **extension field** of a given prime field.

We therefore have to understand field extensions. First note that the field \mathbb{F}' is called an **extension** of a field \mathbb{F} if \mathbb{F} is a subfield of \mathbb{F}' , that is, \mathbb{F} is a subset of \mathbb{F}' and restricting the addition and multiplication laws of \mathbb{F}' to the subset \mathbb{F} recovers the appropriate laws of \mathbb{F} .

Now it can be shown that whenever $p \in \mathbb{P}$ is a prime and $m \in \mathbb{N}$ a natural number, then there is a field \mathbb{F}_{p^m} with characteristic p and p^m elements such that \mathbb{F}_{p^m} is an extension field of the prime field \mathbb{F}_p .

Similarly to the way prime fields \mathbb{F}_p are generated by starting with the ring of integers and then dividing by a prime number p and keeping the remainder, prime field extensions \mathbb{F}_{p^m} are generated by starting with the ring $\mathbb{F}_p[x]$ of polynomials and then dividing them by an irreducible polynomial of degree m and keeping the remainder.

To be more precise, let $P \in F_p[x]$ be an irreducible polynomial of degree m with coefficients from the given prime field \mathbb{F}_p . Then the underlying set \mathbb{F}_{p^m} of the extension field is given by the set of all polynomials with a degree less then m:

$$\mathbb{F}_{p^m} := \{ a_{m-1} x^{m-1} + a_{k-2} x^{k-2} + \ldots + a_1 x + a_0 \mid a_i \in \mathbb{F}_p \}$$
(4.29)

write paragraph on exponentiation

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group pairings

This can be shown to be the set of all remainders when dividing any polynomial $Q \in \mathbb{F}_p[x]$ by P, consequently, elements of the extension field are polynomials of degree less than m. This is analogous to how \mathbb{F}_p is the set of all remainders when dividing integers by p.

Addition is inherited from $\mathbb{F}_p[x]$, which means that addition on \mathbb{F}_{p^m} is defined like normal addition of polynomials:

$$+: \mathbb{F}_{p^m} \times \mathbb{F}_{p^m} \to \mathbb{F}_{p^m}, (\sum_{j=0}^m a_j x^j, \sum_{j=0}^m b_j x^j) \mapsto \sum_{j=0}^m (a_j + b_j) x^j$$
 (4.30)

We can see that the neutral element is (the polynomial) 0, and that the additive inverse is given by the polynomial with all negative coefficients.

Multiplication is inherited from $\mathbb{F}_p[x]$, too, but we have to divide the result by our modulus polynomial P whenever the degree of the resulting polynomial is equal or greater to m:

$$\cdot \mathbb{F}_{p^m} \times \mathbb{F}_{p^m} \to \mathbb{F}_{p^m} , \left(\sum_{j=0}^m a_j x^j, \sum_{j=0}^m b_j x^j \right) \mapsto \left(\sum_{n=0}^{2m} \sum_{i=0}^n a_i b_{n-i} x^n \right) \bmod P \tag{4.31}$$

We can see that the neutral element is (the polynomial) 1. It is, however, not obvious from this definition how the multiplicative inverse looks.

We can easily see from the definition of \mathbb{F}_{p^m} that the field is of characteristic p, since the multiplicative neutral element 1 is equivalent to the multiplicative element 1 from the underlying prime field, and hence $\sum_{j=0}^p 1 = 0$. Moreover, \mathbb{F}_{p^m} is finite and contains p^m many elements, since elements are polynomials of degree < m, and every coefficient a_j can have a p number of different values. In addition, we see that the prime field \mathbb{F}_p is a subfield of \mathbb{F}_{p^m} that occurs when we restrict the elements of \mathbb{F}_p to polynomials of degree zero.

One key point is that the construction of \mathbb{F}_{p^m} depends on the choice of an irreducible polynomial, and, in fact, different choices will give different multiplication tables, since the remainders from dividing a product by P will be different.

It can, however, be shown that the fields for different choices of P are **isomorphic**, which means that there is a one-to-one correspondence between all of them. Consequently, from an abstract point of view, they are the same thing. From an implementations point of view, however, some choices are preferable to others because they allow for faster computations.

To summarize, we have seen that when a prime field \mathbb{F}_p is given, any field \mathbb{F}_{p^m} constructed in the above manner is a field extension of \mathbb{F}_p . To be more general, a field $\mathbb{F}_{p^{m_2}}$ is a field extension of a field $\mathbb{F}_{p^{m_1}}$, if and only if m_1 divides m_2 . From this, we can deduce that, for any given fixed prime number, there are nested sequences of fields whenever the power m_j divides the power m_{j+1} , such that $\mathbb{F}_{p^{m_j}}$ is a subfield of $\mathbb{F}_{p^{m_{j+1}}}$:

$$\mathbb{F}_p \subset \mathbb{F}_{p^{m_1}} \subset \dots \subset \mathbb{F}_{p^{m_k}} \tag{4.32}$$

To get a more intuitive picture of this, we construct an extension field of the prime field \mathbb{F}_3 in the following example, and we can see how \mathbb{F}_3 sits inside that extension field.

Example 63 (The Extension field \mathbb{F}_{3^2}). In (XXX) we have constructed the prime field \mathbb{F}_3 . In this example, we apply the definition of a field extension (page 63) to construct \mathbb{F}_{3^2} . We start by choosing an irreducible polynomial of degree 2 with coefficients in \mathbb{F}_3 . We try $P(t) = t^2 + 1$. Possibly the fastest way to show that P is indeed irreducible is to just insert all elements from

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 \mathbb{F}_3 to see if the result is ever zero. We compute as follows:

$$P(0) = 0^{2} + 1 = 1$$

$$P(1) = 1^{2} + 1 = 2$$

$$P(2) = 2^{2} + 1 = 1 + 1 = 2$$

This implies that P is irreducible. The set \mathbb{F}_{3^2} contains all polynomials of degrees lower than two, with coefficients in \mathbb{F}_3 , which are precisely as listed below:

$$\mathbb{F}_{3^2} = \{0, 1, 2, t, t+1, t+2, 2t, 2t+1, 2t+2\}$$

As expected, our extension field contains 9 elements. Addition is defined as addition of polynomials; for example (t+2)+(2t+2)=(1+2)t+(2+2)=1. Doing this computation for all elements gives the following addition table

+	0	1	2	t	t+1	t+2	2t	2t+1	2t+2
0	0	1	2	t	t+1	t+2	2t	2t+1	2t+2
1	1	2	0	t+1	t+2	t	2t+1	2t+2	2t
2	2	0	1	r+2	t	t+1	2t+2	2t	2t+1
t	t	t+1	t+2	2t	2t+1	2t+2	0	1	2
t+1	t+1	t+2	t	2t+1	2t+2	2t	1	2	0
t+2	t+2	t	t+1	2t+2	2t	2t+1	2	0	1
2t	2t	2t+1	2t+2	0	1	2	t	t+1	t+2
2t+1	2t+1	2t+2	2t	1	2	0	t+1	t+2	t
2t+2	2t+2	2t	2t+1	2	0	1	t+2	t	t+1

As we can see, the group $(\mathbb{F}_3,+)$ is a subgroup of the group $(\mathbb{F}_{3^2},+)$, obtained by only considering the first three rows and columns of this table.

As it was the case in previous examples, we can use the table to deduce the negative of any element from \mathbb{F}_{3^2} . For example, in \mathbb{F}_{3^2} we have -(2t+1)=t+2, since (2t+1)+(t+2)=0

Multiplication needs a bit more computation, as we first have to multiply the polynomials, and whenever the result has a degree ≥ 2 , we have to divide it by P and keep the remainder. To see how this works, let us compute the product of t+2 and 2t+2 in \mathbb{F}_{3^2} :

$$(t+2) \cdot (2t+2) = (2t^2 + 2t + t + 1) \mod (t^2 + 1)$$

$$= (2t^2 + 1) \mod (t^2 + 1) \qquad #2t^2 + 1 : t^2 + 1 = 2 + \frac{2}{t^2 + 1}$$

$$= 2$$

This means that the product of t+2 and 2t+2 in \mathbb{F}_{3^2} is 2. Performing this computation for all elements gives the following multiplication table:

•	0	1	2	t	t+1	t+2	2t	2t+1	2t+2
0	0	0	0	0	0	0	0	0	0
1	0	1	2	t	t+1	t+2	2t	2t+1	2t+2
2	0	2	1	2t	2t+2	2t+1	t	t+2	t+1
t	0	t	2t	2	t+2	2t+2	1	t+1	2t+1
								2	
t+2	0	t+2	2t+1	2t+2	1	t	t+1	2t	2
2t	0	2t	t	1	2t+1	t+1	2	2t+2	t+2
2t+1	0	2t+1	t+2	t+1	2	2t	2t+2	t	1
2t+2	0	2t+2	t+1	2t+1	t	2	t+2	1	2t

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As it was the case in previous examples, we can use the table to deduce the multiplicative inverse of any non-zero element from \mathbb{F}_{3^2} . For example, in \mathbb{F}_{3^2} we have $(2t+1)^{-1}=2t+2$, since $(2t+1)\cdot(2t+2)=1$.

From the multiplication table, we can also see that the only quadratic residues in \mathbb{F}_{3^2} are from the set $\{0,1,2,t,2t\}$, with $\sqrt{0}=\{0\}$, $\sqrt{1}=\{1,2\}$, $\sqrt{2}=\{t,2t\}$, $\sqrt{t}=\{t+2,2t+1\}$ and $\sqrt{2t}=\{t+1,2t+2\}$.

Since \mathbb{F}_{3^2} is a field, we can solve equations as we would for other fields, (such as rational numbers). To see that, let us find all $x \in \mathbb{F}_{3^2}$ that solve the quadratic equation $(t+1)(x^2+(2t+2))=2$. We compute as follows:

Computations in extension fields are arguably on the edge of what can reasonably be done with pen and paper. Fortunately, Sage provides us with a simple way to do the computations.

```
sage: Z3 = GF(3) # prime field
                                                                               192
2424
    sage: Z3t.<t> = Z3[] # polynomials over Z3
                                                                               193
2425
    sage: P = Z3t(t^2+1)
                                                                               194
2426
    sage: P.is_irreducible()
                                                                               195
2427
    True
2428
                                                                               196
    sage: F3_2.<t> = GF(3^2, name='t', modulus=P)
                                                                               197
    sage: F3_2
                                                                               198
2430
    Finite Field in t of size 3<sup>2</sup>
                                                                               199
2431
    sage: F3_2(t+2)*F3_2(2*t+2) == F3_2(2)
                                                                               200
2432
                                                                               201
2433
    sage: F3_2(2*t+2)^(-1) # multiplicative inverse
                                                                               202
2434
    2*t + 1
                                                                               203
2435
    sage: # verify our solution to (t+1)(x^2 + (2t+2)) = 2
                                                                               204
2436
    sage: F3_2(t+1)*(F3_2(t)**2 + F3_2(2*t+2)) == F3_2(2)
                                                                               205
2437
                                                                               206
2438
    sage: F3_2(t+1)*(F3_2(2*t)**2 + F3_2(2*t+2)) == F3_2(2)
                                                                               207
2439
    True
                                                                               208
2440
```

Exercise 32. Consider the extension field \mathbb{F}_{3^2} from the previous example and find all pairs of elements $(x,y) \in \mathbb{F}_{3^2}$, for which the following equation holds:

$$y^2 = x^3 + 4 \tag{4.33}$$

Exercise 33. Show that the polynomial $P = x^3 + x + 1$ from $\mathbb{F}_5[x]$ is irreducible. Then consider the extension field \mathbb{F}_{5^3} defined relative to P. Compute the multiplicative inverse of $(2t^2 + 4) \in$

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 \mathbb{F}_{5^3} using the extended Euclidean algorithm. Then find all $x \in \mathbb{F}_{5^3}$ that solve the following 2445 equation:

$$(2t^2+4)(x-(t^2+4t+2)) = (2t+3)$$
(4.34)

Hashing into extension fields On page 63, we have seen how to hash into prime fields. As check elements of extension fields can be seen as polynomials over prime fields, hashing into extension fields is therefore possible if every coefficient of the polynomial is hashed independently.

Projective Planes 4.4

Projective planes are particular geometric constructs defined over a given field. In a sense, projective planes extend the concept of the ordinary Euclidean plane by including "points at infinity."

Such an inclusion of infinity points makes projective planes particularly useful in the description of elliptic curves, as the description of such a curve in an ordinary plane needs an additional symbol "the point at infinity" to give the set of points on the curve the structure of a group. Translating the curve into projective geometry includes this "point at infinity" more naturally into the set of all points on a projective plane.

To understand the idea of constructing of projective planes, note that in an ordinary Euclidean plane, two lines either intersect in a single point or are parallel. In the latter case, both lines are either the same, that is, they intersect in all points, or do not intersect at all. A projective plane can then be thought of as an ordinary plane, but equipped with additional "point at infinity" such that two different lines always intersect in a single point. Parallel lines intersect "at infinity".

To be more precise, let \mathbb{F} be a field, $\mathbb{F}^3 := \mathbb{F} \times \mathbb{F} \times F$ the set of all three tuples over \mathbb{F} and $x \in \mathbb{F}^3$ with x = (X, Y, Z). Then there is exactly one *line* in \mathbb{F}^3 that intersects both (0,0,0) and x. This line is given as follows:

$$[X:Y:Z] := \{(k \cdot X, k \cdot Y, k \cdot Z) \mid k \in \mathbb{F}\}$$

$$(4.35)$$

A **point** in the **projective plane** over \mathbb{F} is defined as such a **line**, and the projective plane is the 2468 set of all such points: 2469

$$\mathbb{FP}^2 := \{ [X : Y : Z] \mid (X, Y, Z) \in \mathbb{F}^3 \text{ with } (X, Y, Z) \neq (0, 0, 0) \}$$
(4.36)

It can be shown that a projective plane over a finite field \mathbb{F}_{p^m} contains $p^{2m} + p^m + 1$ number of 2470 elements. 2471

To understand why [X:Y:Z] is called a line, consider the situation where the underlying field \mathbb{F} is the set of real numbers \mathbb{R} . In this case, \mathbb{R}^3 can be seen as the three-dimensional space, and [X:Y:Z] is an ordinary line in this 3-dimensional space that intersects zero and the point with coordinates X, Y and Z.

The key observation here is that points in the projective plane are lines in the 3-dimensional space \mathbb{F}^3 ,]. Additionally, for finite fields, the terms **space** and **line** share very little visual similarity with their counterparts over the set of real numbers.

It follows from this that points $[X:Y:Z] \in \mathbb{FP}^2$ are not simply described by fixed coordinates (X,Y,Z), but by sets of coordinates, where two different coordinates (X_1,Y_1,Z_1) and (X_2, Y_2, Z_2) describe the same point if and only if there is some field element k such that $(X_1,Y_1,Z_1)=(k\cdot X_2,k\cdot Y_2,k\cdot Z_2)$. Points [X:Y:Z] are called **projective coordinates**.

Notation and Symbols 6 (Projective coordinates). Projective coordinates of the form [X:Y:1] are descriptions of so-called **affine points**. Projective coordinates of the form [X:Y:0] are descriptions of so-called **points at infinity**. In particular, the projective coordinate [1:0:0] describes the so-called **line at infinity**.

Example 64. Consider the field \mathbb{F}_3 from example XXX. As this field only contains three elements, it does not take too much effort to construct its associated projective plane $\mathbb{F}_3\mathbb{P}^2$, as we know that it only contains 13 elements.

add reference

To find $\mathbb{F}_3\mathbb{P}^2$, we have to compute the set of all lines in $\mathbb{F}_3 \times \mathbb{F}_3 \times \mathbb{F}_3$ that intersect (0,0,0). Since those lines are parameterized by tuples (x_1,x_2,x_3) , we compute as follows:

```
[0:0:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,0,1), (0,0,2)\}
[0:0:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,0,2), (0,0,1)\} = [0:0:1]
[0:1:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,1,0), (0,2,0)\}
[0:1:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,1,1), (0,2,2)\}
[0:1:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,1,2), (0,2,1)\}
[0:2:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,2,0), (0,1,0)\} = [0:1:0]
[0:2:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,2,1), (0,1,2)\} = [0:1:2]
[0:2:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(0,2,2), (0,1,1)\} = [0:1:1]
[1:0:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,0,0), (2,0,0)\}
[1:0:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,0,1), (2,0,2)\}
[1:0:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,0,2), (2,0,1)\}
[1:1:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,1,0), (2,2,0)\}
[1:1:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,1,1), (2,2,2)\}
[1:1:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,1,2), (2,2,1)\}
[1:2:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,2,0), (2,1,0)\}
[1:2:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,2,1), (2,1,2)\}
[1:2:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(1,2,2), (2,1,1)\}
[2:0:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,0,0), (1,0,0)\} = [1:0:0]
[2:0:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,0,1), (1,0,2)\} = [1:0:2]
[2:0:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,0,2), (1,0,1)\} = [1:0:1]
[2:1:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,1,0), (1,2,0)\} = [1:2:0]
[2:1:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,1,1), (1,2,2)\} = [1:2:2]
[2:1:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,1,2), (1,2,1)\} = [1:2:1]
[2:2:0] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,2,0), (1,1,0)\} = [1:1:0]
[2:2:1] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,2,1), (1,1,2)\} = [1:1:2]
[2:2:2] = \{(k \cdot x_1, k \cdot x_2, k \cdot x_3) \mid k \in \mathbb{F}_3\} = \{(2,2,2), (1,1,1)\} = [1:1:1]
```

These lines define the 13 points in the projective plane $\mathbb{F}_3\mathbb{P}$:

```
\mathbb{F}_{3}\mathbb{P} = \{[0:0:1], [0:1:0], [0:1:1], [0:1:2], [1:0:0], [1:0:1], \\ [1:0:2], [1:1:0], [1:1:1], [1:1:2], [1:2:0], [1:2:1], [1:2:2]\}
```

This projective plane contains 9 affine points, three points at infinity and one line at infinity.

To understand the ambiguity in projective coordinates a bit better, let us consider the point [1:2:2]. As this point in the projective plane is a line in \mathbb{F}_3^3 , it has the projective coordinates (1,2,2) as well as (2,1,1), since the former coordinate gives the latter when multiplied in \mathbb{F}_3 by the factor 2. In addition, note that, for the same reasons, the points [1:2:2] and [2:1:1] are the same, since their underlying sets are equal.

Exercise 34. Construct the so-called **Fano plane**, that is, the projective plane over the finite field \mathbb{F}_2 .

Chapter 5

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Elliptic Curves

Generally speaking, elliptic curves are "curves" defined in geometric planes like the Euclidean or the projective plane over some given field. One of the key features of elliptic curves over finite fields from the point of view of cryptography is that their set of points has a group law such that the resulting group is finite and cyclic, and it is believed that the discrete logarithm problem on these groups is hard.

A special class of elliptic curves are so-called **pairing-friendly curves**, which have a notation of a group pairing as defined in XXX. This pairing has cryptographically advantageous properties. Those curve are useful in the development of SNARKs, since they allow to compute so-called R1CS-satisfiability "in the exponent" MIRCO: (THIS HAS TO BE REWRITTEN WITH WAY MORE DETAIL)

In this chapter, we introduce epileptic curves as they are used in pairing-based approaches to the construction of SNARKs. The elliptic curves we consider are all defined over prime fields or prime field extensions and the reader should be familiar with the contend of the previous section on those fields.

In its most generality elliptic curves are defined as a smooth projective curve of genus 1 defined over some field $\mathbb F$ with a distinguished $\mathbb F$ -rational point, but this definition is not very useful for the introductory character of this book. We will therefore look at 3 more practical definitions in the following sections, by introducing Weierstraß, Montgomery and Edwards curves. All of them are widely used in cryptography, and understanding them is crucial to being able to follow the rest of this book.

TODO: Elliptic Curve asymmetric cryptography examples. Private key, generator, public key.

add reference

maybe remove this sentence?

5.1 Elliptic Curve Arithmetics

5.1.1 Short Weierstraß Curves

In this section, we introduce **short Weierstraß** curves, which are the most general types of curves over finite fields of characteristic greater than 3.

We start with their representation in affine space. This representation has the advantage that affine points correspond to pairs of numbers, which makes it more accessible for beginners. However, it has the disadvantage that a special "point at infinity", that is not a point on the curve, is necessary to describe the group structure. We introduce the elliptic curve group law and describe elliptic curve scalar multiplication, which is an instantiation of the exponential map from general cyclic groups.

Then we look at the projective representation of short Weierstraß curves. This has the advantage that no special symbol is necessary to represent the point at infinity but comes with

affine space

the drawback that projective points are classes of numbers, which might be a bit unusual for a beginner.

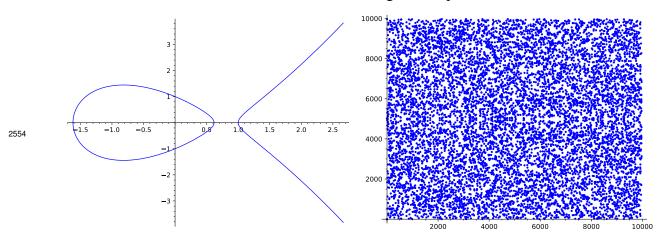
We finish this section with an explicit equivalence that transforms affine representations into projective ones and vice versa.

Affine short Weierstraß form Probably the least abstract and most straight-forward way to introduce elliptic curves for non-mathematicians and beginners is the so-called affine representation of a short Weierstraß curve. To see what this is, let \mathbb{F} be a finite field of order q and $a, b \in \mathbb{F}$ two field elements such that $4a^3 + 27b^2 \mod q \neq 0$. Then a **short Weierstraß elliptic curve** $E(\mathbb{F})$ over \mathbb{F} in its affine representation is the set of all pairs of field elements $(x,y) \in \mathbb{F} \times \mathbb{F}$ that satisfy the short Weierstraß cubic equation $y^2 = x^3 + a \cdot x + b$, together with a distinguished symbol \mathcal{O} , called the **point at infinity**:

$$E(\mathbb{F}) = \{(x, y) \in \mathbb{F} \times \mathbb{F} \mid y^2 = x^3 + a \cdot x + b\} \bigcup \{\mathscr{O}\}$$
 (5.1)

Notation and Symbols 7. In the literature, the set $E(\mathbb{F})$, which includes the symbol \mathcal{O} , is often called the set of **rational points** of the elliptic curve, in which case the curve itself is usually written as E/\mathbb{F} . However, in what follows, we will frequently identify an elliptic curve with its set of rational points and therefore use the notation $E(\mathbb{F})$ instead. This is possible in our case, since we only the group structure of the curve in consideration is relevant for us.

The term "curve" is used here because, in the ordinary 2 dimensional plane \mathbb{R}^2 , the set of all points (x,y) that satisfy $y^2 = x^3 + a \cdot x + b$ looks like a curve. We should note however that visualizing elliptic curves over finite fields as "curves" has its limitations, and we will therefore not stress the geometric picture too much, but focus on the computational properties instead. To understand the visual difference, consider the following two elliptic curves:



Both elliptic curves are defined by the same short Weierstraß equation $y^2 = x^3 - 2x + 1$, but the first curve is defined in the real affine plane \mathbb{R}^2 , that is, the pair (x,y) contains real numbers, while the second one is defined in the affine plane \mathbb{F}^2_{9973} , which means that both x and y are from the prime field \mathbb{F}_{9973} . Every blue dot represents a pair (x,y), that is a solution to $y^2 = x^3 - 2x + 1$. As we can see, the second curve hardly looks like a geometric structure one would naturally call a curve. This shows that our geometric intuitions from \mathbb{R}^2 are obfuscated in curves over finite fields.

The identity $6 \cdot (4a^3 + 27b^2) \mod q \neq 0$ ensures that the curve is non-singular, which basically means that the curve has no cusps or self-intersections.

self-intersection

Throughout this book, the reader is advised to do as many computations in a pen-and-paper fashion as possible, as this is helps getting a deeper understanding of the details. However, when dealing with elliptic curves, computations can quickly become cumbersome and tedious, and one might get lost in the details. Fortunately, Sage is very helpful in dealing with elliptic curves. This book to introduces the reader to the great elliptic curve capabilities of Sage. The following snippet shows a way to define elliptic curves and work with them in Sage:

```
sage: F5 = GF(5) # define the base field
                                                                              209
2570
    sage: a = F5(2) # parameter a
                                                                              210
2571
    sage: b = F5(4) # parameter b
                                                                              211
2572
    sage: # check non-sigularity
                                                                              212
2573
    sage: F5(6)*(F5(4)*a^3+F5(27)*b^2) != F5(0)
2574
                                                                              213
    True
                                                                              214
2575
    sage: # short Weierstrass curve
                                                                              215
2576
    sage: E = EllipticCurve(F5, [a,b]) # y^2 == x^3 + ax + b
                                                                              216
2577
    sage: P = E(0,2) \# 2^2 == 0^3 + 2*0 + 4
                                                                              217
2578
    sage: P.xy() # affine coordinates
                                                                              218
2579
    (0, 2)
                                                                              219
2580
    sage: INF = E(0) # point at infinity
                                                                              220
2581
                  # point at infinity has no affine coordinates
                                                                              221
2582
                INF.xy()
                                                                              222
2583
    ....: except ZeroDivisionError:
                                                                              223
2584
                                                                              224
                pass
2585
    sage: P = E.plot() # create a plotted version
                                                                              225
2586
```

The following three examples give a more practical understanding of what an elliptic curve is and how we can compute it. The reader is advised to read them carefully, and ideally, to also carry out the computation themselves. We will repeatedly build on these examples in this chapter, and use the second example throughout this book.

Example 65. To provide the reader with an example of a small elliptic curve where all computation can be done with pen and paper, consider the prime field \mathbb{F}_5 from example 59 (page 61). quite familiar to readers who had worked through the examples and exercises in the previous chapter.

check reference

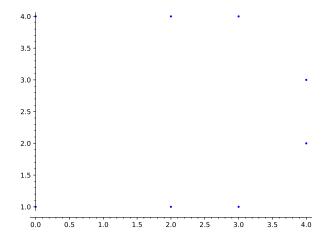
To define an elliptic curve over \mathbb{F}_5 , we have to choose to numbers a and b from that field. Assuming we choose a=1 and b=1 then $4a^3+27b^2\equiv 1\pmod 5$ from which follows that the corresponding elliptic curve $E_1(\mathbb{F}_5)$ is given by the set of all pairs (x,y) from \mathbb{F}_5 that satisfy the equation $y^2=x^3+x+1$, together with the special symbol \mathscr{O} , which represents the "point at infinity".

To get a better understand of that curve, observer that if we choose arbitrarily the pair (x,y) = (1,1), we see that $1^2 \neq 1^3 + 1 + 1$ and hence (1,1) is not an element of the curve $E_1(\mathbb{F}_5)$. On the other hand choosing for example (x,y) = (2,1) gives $1^2 = 2^3 + 2 + 1$ and hence the pair (2,1) is an element of $E_1(\mathbb{F}_5)$ (Remember that all computations are done in modulo 5 arithmetics).

Now since the set $\mathbb{F}_5 \times \mathbb{F}_5$ of all pairs (x,y) from \mathbb{F}_5 contains only $5 \cdot 5 = 25$ pairs, we can compute the curve, by just inserting every possible pair (x,y) into the short Weierstraß equation $y^2 = x^3 + x + 1$. If the equation holds, the pair is a curve point, if not that means that the point is not on the curve. Combining the result of this computation with the point at infinity gives the curve as follows:

$$E_1(\mathbb{F}_5) = \{ \mathscr{O}, (0,1), (2,1), (3,1), (4,2), (4,3), (0,4), (2,4), (3,4) \}$$

This means that our elliptic curve is a set of 9 elements, 8 of which are pairs of numbers and one special symbol \mathcal{O} . Visualizing E1 gives the following plot:



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In the development of SNARKs, it is sometimes necessary to do elliptic curve cryptography "in a circuit", which basically means that the elliptic curves need to be implemented in a certain SNARK-friendly way. We will look at what this means in chapter 7. To be able to do this efficiently, it is desirable to have curves with special properties. The following example is a pen-and-paper version of such a curve, called **Baby-jubjub**, which resembles cryptographically secure curves extensively used in real-world SNARKs. The interested reader is advised to study this example carefully, as we will use it and build on it in various places throughout the book. I feel like a lot of people won't get the Lewis Carroll reference unless we make it more explicit

check reference

jubjub

Example 66 (Pen-JubJub). Consider the prime field \mathbb{F}_{13} from exercise 4.3 (page 63. If we choose a=8 and b=8, then $4a^3+27b^2\equiv 6\pmod{13}$ and the corresponding elliptic curve is given by all pairs (x,y) from \mathbb{F}_{13} such that $y^2=x^3+8x+8$ holds. We call this curve the **Pen-JubJub** curve, or PJJ_13 for short.

check reference

Now, since the set $\mathbb{F}_{13} \times \mathbb{F}_{13}$ of all pairs (x,y) from \mathbb{F}_{13} contains only $13 \cdot 13 = 169$ pairs, we can compute the curve by just inserting every possible pair (x,y) into the short Weierstraß equation $y^2 = x^3 + 8x + 8$. We get the following result:

$$PJJ_13 = \{ \mathscr{O}, (1,2), (1,11), (4,0), (5,2), (5,11), (6,5), (6,8), (7,2), (7,11), (8,5), (8,8), \\ (9,4), (9,9), (10,3), (10,10), (11,6), (11,7), (12,5), (12,8) \}$$
 (5.2)

As we can see, the curve consists of 20 points; 19 points from the affine plane and the point at infinity. To get a visual impression of the PJJ_13 curve, we might plot all of its points (except the point at infinity) in the $\mathbb{F}_{13} \times \mathbb{F}_{13}$ affine plane. We get the following plot:

affine plane

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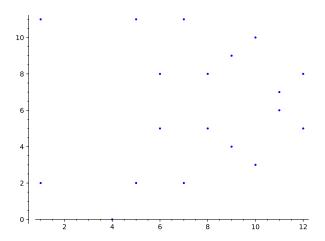
2639

2640

2641

2642

2643



As we will see in what follows, this curve is rather special, as it is possible to represent it in two alternative forms called the **Montgomery** and the **twisted Edwards form** (See sections 5.1.2, and 5.1.3, respectively).

check reference

Now that we have seen two pen-and-paper friendly elliptic curves, let us look at a curve, that is used in actual cryptography. Cryptographically secure elliptic curves are not **qualitatively** different from the curves we looked at so far, but the prime number modulus of their prime field is much larger. Typical examples use prime numbers that have binary representations in the magnitude of more than double the size of the desired security level. If, for example, a security of 128 bits is desired, a prime modulus of binary size \geq 256 is chosen. The following example provides such a curve.

check

reference

Example 67 (Bitcoin's Secp256k1 curve). To give an example of a real-world, cryptographically secure curve, let us look at curve Secp256k1, which is famous for being used in the public key cryptography of Bitcoin. The prime field \mathbb{F}_p of Secp256k1 is defined by the following prime number:

p = 115792089237316195423570985008687907853269984665640564039457584007908834671663

The binary representation of this number needs 256 bits, which implies that the prime field \mathbb{F}_p contains approximately 2^{256} many elements, which is considered quite large. To get a better impression of how large the base field is, consider that the number 2^{256} is approximately in the same order of magnitude as the estimated number of atoms in the observable universe.

The curve Secp256k1 is defined by the parameters $a, b \in \mathbb{F}_p$ with a = 0 and b = 7. Since $4 \cdot 0^3 + 27 \cdot 7^2 \mod p = 1323$, those parameters indeed define an elliptic curve given as follows:

$$Secp256k1 = \{(x, y) \in \mathbb{F}_p \times \mathbb{F}_p \mid y^2 = x^3 + 7 \}$$

Clearly, the Secp256k1 curve is too large to do computations by hand, since it can be shown that the number of its elements is a prime number r that also has a binary representation of 256 bits:

r = 115792089237316195423570985008687907852837564279074904382605163141518161494337

Cryptographically secure elliptic curves are therefore not useful in pen-and-paper computations. Fortunately, Sage handles large curves efficiently:

sage: p = 1157920892373161954235709850086879078532699846656405 226
64039457584007908834671663

```
sage: # Hexadecimal representation
                                                                       227
2644
    sage: p.str(16)
                                                                       228
2645
    229
2646
2647
    sage: p.is_prime()
                                                                       230
2648
    True
                                                                       231
2649
    sage: p.nbits()
                                                                       232
2650
    256
                                                                       233
2651
    sage: Fp = GF(p)
                                                                       234
2652
    sage: Secp256k1 = EllipticCurve(Fp, [0, 7])
2653
                                                                       235
    sage: r = Secp256k1.order() # number of elements
                                                                       236
2654
    sage: r.str(16)
                                                                       237
2655
    ffffffffffffffffffffffffffffbaaedce6af48a03bbfd25e8cd03641
                                                                       238
2656
       41
2657
    sage: r.is_prime()
                                                                       239
2658
    True
                                                                       240
2659
    sage: r.nbits()
                                                                       241
2660
    256
                                                                       242
2661
```

Exercise 35. Look up the definition of curve BLS12-381, implement it in Sage and compute its order.

Affine compressed representation As we have seen in example 67, cryptographically secure elliptic curves are defined over large prime fields, where elements of those fields typically need more than 255 bits of storage on a computer. Since elliptic curve points consist of pairs of those field elements, they need double that amount of storage.

check reference

However, we can reduce the amount of space needed to represent a curve point by using a technique called **point compression**. Note that, up to a sign, the y coordinate of a curve point can be computed from the x coordinate by simply inserting x into the Weierstraß equation and then computing the roots of the result. This gives two results, and it means that we can represent a curve point in **compressed form** by simply storing the x coordinate together with a single sign bit only, the latter of which deterministically decides which of the two roots to choose. One convention could be to always choose the root closer to 0 when the sign bit is 0, and the root closer to the order of \mathbb{F} when the sign bit is 1. In case the y coordinate is zero, both sign bits give the same result.

more explanation of what

sign

Example 68 (Pen-jubjub). To understand the concept of compressed curve points a bit better, consider the PJJ_13 curve from example 66 again. Since this curve is defined over the prime field \mathbb{F}_{13} , and numbers between 0 and 13 need approximately 4 bits to be represented, each PJJ_13 point on this curve needs 8 bits of storage in uncompressed form. The following set represents the uncompressed form of the points on this curve:

check reference

the sign is

$$PJJ_13 = \{ \mathscr{O}, (1,2), (1,11), (4,0), (5,2), (5,11), (6,5), (6,8), (7,2), (7,11), (8,5), (8,8), (9,4), (9,9), (10,3), (10,10), (11,6), (11,7), (12,5), (12,8) \}$$

Using the technique of point compression, we can reduce the bits needed to represent the points on this curve to 5 per point. To achieve this, we can replace the y coordinate in each (x, y) pair by a sign bit indicating whether or not y is closer to 0 or to 13. As a result y values in the range [0, ..., 6] will have the sign bit 0, while y-values in the range [7, ..., 12] will have the sign bit 1.

Applying this to the points in *PJJ_13* gives the compressed representation as follows:

```
PJJ\_13 = \{ \mathscr{O}, (1,0), (1,1), (4,0), (5,0), (5,1), (6,0), (6,1), (7,0), (7,1), \\ (8,0), (8,1), (9,0), (9,1), (10,0), (10,1), (11,0), (11,1), (12,0), (12,1) \}
```

Note that the numbers 7, ..., 12 are the negatives (additive inverses) of the numbers 1, ..., 6 in modular 13 arithmetics and that -0 = 0. Calling the compression bit a "sign bit" therefore makes sense.

To recover the uncompressed counterpart of, say, the compressed point (5,1), we insert the x coordinate 5 into the Weierstraß equation and get $y^2 = 5^3 + 8 \cdot 5 + 8 = 4$. As expected, 4 is a quadratic residue in \mathbb{F}_{13} with roots $\sqrt{4} = \{2,11\}$. Since the sign bit of the point is 1, we have to choose the root closer to the modulus 13, which is 11. The uncompressed point is therefore (5,11).

S: I don't follow this at all

Looking at the previous examples, the compression rate does not look very impressive. However, looking at the real-life example of the Secp256k1 curve shows that compression is has significant practical advantages.

Example 69. Consider the Secp256k1 curve from example 67 again. The following code invokes Sage to generate a random affine curve point, then applies our compression method to it:

check reference

```
sage: P = Secp256k1.random_point().xy()
                                                                              243
2691
    sage: P
                                                                              244
2692
    (5732745559092928700275495328195703081931555862512446945836228
                                                                              245
2693
       5630887028852436, 24242609999426606897142811967939071817174
2694
       686615886596221090801834998454950146)
2695
    sage: # uncompressed affine point size
                                                                              246
    sage: ZZ(P[0]).nbits()+ZZ(P[1]).nbits()
                                                                              247
2697
    509
                                                                              248
2698
    sage: # compute the compression
                                                                              249
2699
    sage: if P[1] > Fp(-1)/Fp(2):
                                                                              250
2700
                PARITY = 1
    . . . . :
                                                                              251
2701
    ....: else:
                                                                              252
2702
                PARITY = 0
2703
                                                                              253
    sage: PCOMPRESSED = [P[0], PARITY]
                                                                              254
2704
    sage: PCOMPRESSED
                                                                              255
2705
    [5732745559092928700275495328195703081931555862512446945836228
                                                                             256
2706
       5630887028852436, 0]
2707
    sage: # compressed affine point size
                                                                              257
2708
    sage: ZZ(PCOMPRESSED[0]).nbits()+ZZ(PCOMPRESSED[1]).nbits()
                                                                              258
2709
                                                                              259
    255
                                                                               add ex-
2711
```

define of how this shows what we claim

planation

Affine group law One of the key properties of an elliptic curve is that it is possible to define a group law on the set of its rational points such that the point at infinity serves as the neutral element and inverses are reflections on the *x*-axis.

The origin of this law can be understood in a geometric picture and is known as the **chord-and-tangent rule**. In the affine representation of a short Weierstraß curve, the rule can be described in the following way:

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Definition 5.1.1.1. Chord-and-tangent rule

- (Point at infinity) We define the point at infinity \mathscr{O} as the neutral element of addition, that is, we define $P + \mathscr{O} = P$ for all points $P \in E(\mathbb{F})$.
- (Point addition) Let $P,Q \in E(\mathbb{F}) \setminus \{ \mathcal{O} \}$ with $P \neq Q$ be two distinct points on an elliptic curve, neither of them the point at infinity. The sum of P and Q is defined as follows: Consider the line l which intersects the curve in P and Q. If l intersects the elliptic curve at a third point R', define the sum $R = P \oplus Q$ of P and Q as the reflection of R' at the x-axis. If the line l does not intersect the curve at a third point, define the sum to be the point at infinity \mathcal{O} . It can be shown that no such chord line will intersect the curve in more than three points, so addition is not ambiguous.
- (Point doubling) Let P∈ E(F)\{Ø} be a point on an elliptic curve, that is not the point at infinity. The sum of P with itself (the doubling of P) is defined as follows:
 Consider the line which is tangential to the elliptic curve at P. If this line intersects the elliptic curve at a second point R', the sum 2P = P + P is the reflection of R' at the x-axis. If it does not intersect the curve at a third, point define the sum to be the point at infinity Ø. It can be shown that no such tangent line will intersect the curve in more than two points, so addition is not ambiguous.

It can be shown that the points of an elliptic curve form a commutative group with respect to the tangent-and-chord rule such that \mathscr{O} acts the neutral element, and the inverse of any element $P \in E(\mathbb{F})$ is the reflection of P on the x-axis.

To translate the geometric description into algebraic equations, first observe that, for any two given curve points $(x_1, y_1), (x_2, y_2) \in E(\mathbb{F})$, it can be shown that the identity $x_1 = x_2$ implies $y_2 = \pm y_1$, which shows that the following rules are a complete description of the affine addition law.

Definition 5.1.1.2. Chord-and-tangent rule: algebraic equations

- (Neutral element) The point at infinity \mathcal{O} is the neutral element.
- (Additive inverse) The additive inverse of \mathscr{O} is \mathscr{O} . For any other curve point $(x,y) \in E(\mathbb{F}) \setminus \{\mathscr{O}\}$, the additive inverse is given by (x,-y).
- (Addition rule) For any two curve points $P,Q \in E(\mathbb{F})$, addition is defined by one of the following three cases:
 - 1. (Adding the neutral element) If $Q = \mathcal{O}$, then the sum is defined as $P \oplus Q = P$.
 - 2. (Adding inverse elements) If P = (x, y) and Q = (x, -y), the sum is defined as $P \oplus Q = \emptyset$.
 - 3. (Adding non-self-inverse equal points) If P = (x, y) and Q = (x, y) with $y \ne 0$, the sum 2P = (x', y') is defined as follows: We only referred to P in the definition of point doubling above so Q seems a bit confusing here even though it's defined as equal to P

$$x' = \left(\frac{3x^2+a}{2y}\right)^2 - 2x$$
 , $y' = \left(\frac{3x^2+a}{2y}\right)^2 (x-x') - y$

4. (Adding non-inverse different points) If $P = (x_1, y_1)$ and $Q = (x_2, y_2)$ such that $x_1 \neq x_2$, the sum R = P + Q with $R = (x_3, y_3)$ is defined as follows:

$$x_3 = \left(\frac{y_2 - y_1}{x_2 - x_1}\right)^2 - x_1 - x_2$$
, $y_3 = \left(\frac{y_2 - y_1}{x_2 - x_1}\right)(x_1 - x_3) - y_1$

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this def. be moved even earlier?

should

chord line

tangential

tangent line

remove

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Note that short Weierstraß curve points P with P = (x, 0) are inverses of themselves, which 2751 implies $2P = \mathcal{O}$ in this case.

Notation and Symbols 8. Let \mathbb{F} be a field and $E(\mathbb{F})$ be an elliptic curve over \mathbb{F} . We write \oplus for the group law on $E(\mathbb{F})$ and $(E(\mathbb{F}), \oplus)$ for the group of rational points.

As we can see, it is very efficient to compute inverses on elliptic curves. However, computing the addition of elliptic curve points in the affine representation needs to consider many cases and involves extensive finite field divisions. As we will see in the next paragraph, this can be simplified in projective coordinates.

where?

To get some practical impression of how the group law on an elliptic curve is computed, let's look at some actual cases:

Example 70. Consider the elliptic curve $E_1(\mathbb{F}_5)$ from example 65 again. As we have seen, the set of rational points contains 9 elements:

check reference

$$E_1(\mathbb{F}_5) = \{ \mathscr{O}, (0,1), (2,1), (3,1), (4,2), (4,3), (0,4), (2,4), (3,4) \}$$

We know that this set defines a group, so we can add any two elements from $E_1(\mathbb{F}_5)$ to get a third element.

To give an example, consider the elements (0,1) and (4,2). Neither of these elements is the neutral element \mathcal{O} , and since, the x coordinate of (0,1) is different from the x coordinate of (4,2), we know that we have to use the chord rule, that is, rule number 4 from definition 5.1.1.2 to compute the sum $(0,1) \oplus (4,2)$:

check reference

$$x_{3} = \left(\frac{y_{2} - y_{1}}{x_{2} - x_{1}}\right)^{2} - x_{1} - x_{2}$$
insert points
$$= \left(\frac{2 - 1}{4 - 0}\right)^{2} - 0 - 4$$
simplify in \mathbb{F}_{5}

$$= \left(\frac{1}{4}\right)^{2} + 1 = 4^{2} + 1 = 1 + 1 = 2$$

$$y_3 = \left(\frac{y_2 - y_1}{x_2 - x_1}\right) (x_1 - x_3) - y_1$$
 # insert points

$$= \left(\frac{2 - 1}{4 - 1}\right) (0 - 2) - 1$$
 # simplify in \mathbb{F}_5

$$= \left(\frac{1}{4}\right) \cdot 3 + 4 = 4 \cdot 3 + 4 = 2 + 4 = 1$$

So, in our elliptic curve $E_1(\mathbb{F}_5)$ we get $(0,1) \oplus (4,2) = (2,1)$, and, indeed, the pair (2,1) is an element of $E_1(\mathbb{F}_5)$ as expected. On the other hand, $(0,1) \oplus (0,4) = \emptyset$, since both points have equal x coordinates and inverse y coordinates, rendering them inverses of each other. Adding the point (4,2) to itself, we have to use the tangent rule, that is, rule 3 from definition 5.1.1.2:

check reference

$$x' = \left(\frac{3x^2 + a}{2y}\right)^2 - 2x$$
insert points
$$= \left(\frac{3 \cdot 4^2 + 1}{2 \cdot 2}\right)^2 - 2 \cdot 4$$

$$= \left(\frac{3 \cdot 1 + 1}{4}\right)^2 + 3 \cdot 4 = \left(\frac{4}{4}\right)^2 + 2 = 1 + 2 = 3$$

$$y' = \left(\frac{3x^2 + a}{2y}\right)^2 (x - x') - y$$
insert points

$$y' = \left(\frac{2y}{2y}\right) (x - x') - y$$

$$= \left(\frac{3 \cdot 4^2 + 1}{2 \cdot 2}\right)^2 (4 - 3) - 2$$

$$= 1 \cdot 1 + 3 = 4$$
insert points
$$= \left(\frac{3 \cdot 4^2 + 1}{2 \cdot 2}\right)^2 (4 - 3) - 2$$

$$= 1 \cdot 1 + 3 = 4$$

So, in our elliptic curve $E_1(\mathbb{F}_5)$, we get the doubling of (4,2), that is, $(4,2) \oplus (4,2) = (3,4)$, 2763 and, indeed the pair (3,4) is an element of $E_1(\mathbb{F}_5)$ as expected. The group $E_1(\mathbb{F}_5)$ has no selfinverse points other than the neutral element \mathcal{O} , since no point has 0 as its y coordinate. We can 2765 invoke Sage to double-check the computations. 2766

```
sage: F5 = GF(5)
                                                                                    260
2767
    sage: E1 = EllipticCurve(F5,[1,1])
                                                                                    261
2768
    sage: INF = E1(0) # point at infinity
                                                                                    262
2769
    sage: P1 = E1(0,1)
                                                                                    263
2770
    sage: P2 = E1(4,2)
                                                                                    264
2771
    sage: P3 = E1(0,4)
                                                                                    265
    sage: R1 = E1(2,1)
                                                                                    266
2773
    sage: R2 = E1(3,4)
                                                                                    267
2774
    sage: R1 == P1+P2
                                                                                    268
2775
    True
                                                                                    269
2776
    sage: INF == P1+P3
                                                                                    270
2777
    True
                                                                                    271
2778
    sage: R2 == P2+P2
                                                                                    272
                                                                                    273
2780
    sage: R2 == 2*P2
                                                                                    274
2781
                                                                                    275
2782
    sage: P3 == P3 + INF
                                                                                    276
2783
    True
                                                                                    277
2784
```

Example 71 (Pen-jubjub). Consider the PJJ_13-curve from example 66 again and recall that its check group of rational points is given as follows:

reference

$$PJJ_{1} = \{ \mathscr{O}, (1,2), (1,11), (4,0), (5,2), (5,11), (6,5), (6,8), (7,2), (7,11), (8,5), (8,8), (9,4), (9,9), (10,3), (10,10), (11,6), (11,7), (12,5), (12,8) \}$$

In contrast to the group from the previous example, this group contains a self-inverse point, 2785 which is different from the neutral element, defined by (4,0). To see what this means, observe 2786 that we cannot add (4,0) to itself using the tangent rule 3 from definition 5.1.1.2, as the y check 2787 coordinate is zero. Instead, we have to use rule 2, since 0 = -0. We therefore get $(4,0) \oplus$ reference $(4,0) = \mathcal{O}$ in PJJ_13 . The point (4,0) is therefore the inverse of itself, as adding it to itself results in the neutral element.

```
sage: F13 = GF(13)
                                                                                 278
2791
    sage: MJJ = EllipticCurve(F13,[8,8])
                                                                                 279
2792
    sage: P = MJJ(4,0)
                                                                                 280
2793
    sage: INF = MJJ(0) # Point at infinity
                                                                                 281
2794
           INF == P+P
    sage:
                                                                                 282
2795
    True
                                                                                 283
2796
    sage: INF == 2*P
                                                                                 284
2797
    True
                                                                                 285
```

Example 72. Consider the Secp256k1 curve from example 67 <u>again</u>. The following code invokes Sage to generate a random affine curve point, then applies our compression method:

check reference

```
sage: P = Secp256k1.random_point()
                                                                            286
2801
    sage: Q = Secp256k1.random_point()
                                                                            287
2802
    sage: INF = Secp256k1(0)
                                                                            288
2803
    sage: R1 = -P
                                                                            289
2804
    sage: R2 = P + Q
                                                                            290
2805
    sage: R3 = Secp256k1.order()*P
                                                                            291
2806
                                                                            292
    sage: P.xy()
2807
    (2437965124411773648884901383952245798298026200193112014924104
                                                                           293
2808
       5920541255603582, 38155318538062562663408568861188374070643
2809
       301057931057692802349663368915027747)
2810
                                                                            294
2811
    sage: Q.xy()
    (6273267811834346524071370277009541823203325405903695727983144
                                                                           295
2812
       7554159754801518, 81206263702504109131546480004400274036228
2813
       732572045186080577817223096074627142)
2814
    sage: (ZZ(R1[0]).str(16), ZZ(R1[1]).str(16))
                                                                            296
2815
    ('35e664c3768462813f30192e327e60c61508d279931cdbc639f3cb11c5b3
                                                                           297
2816
       157e', 'aba4dae1f8c83f0ac955259cd78622327b9f107d82937463dd8
2817
       cded0c012750c')
2818
    sage: R2.xy()
                                                                            298
2819
    (8315162076242884051827668971975027473477042355284820491860209
                                                                            299
2820
       945466147353499, 128083043736478847072934266448265932843478
2821
       45733596286872839204967881615931190)
2822
    sage: R3 == INF
                                                                            300
2823
                                                                            301
2824
    sage: P[1]+R1[1] == Fp(0) \# -(x,y) = (x,-y)
                                                                            302
2825
    True
                                                                            303
2826
```

2827 Exercise 36. Consider the PJJ_13-curve from example 66.

check reference

- 1. Compute the inverse of (10,10), \mathcal{O} , (4,0) and (1,2).
- 2829 2. Compute the expression 3*(1,11)-(9,9).

2828

- 2830 3. Solve the equation x + 2(9,4) = (5,2) for some $x \in PJJ_13$
- 4. Solve the equation $x \cdot (7,11) = (8,5)$ for $x \in \mathbb{Z}$

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Scalar multiplication As we have seen in the previous section, elliptic curves $E(\mathbb{F})$ have the structure of a commutative group associated to them. Moreover, It can moreover be shown that this group is finite and cyclic whenever the field is finite.

To understand elliptic curve scalar multiplication, recall from page 43 that every finite cyclic check group of order q has a generator g and an associated exponential map $g^{(\cdot)}: \mathbb{Z}_q \to \mathbb{G}$, where g^n is the *n*-fold product of *g* with itself.

reference

Elliptic curve scalar multiplication is the exponential map written in additive notation. To be more precise, let \mathbb{F} be a finite field, $E(\mathbb{F})$ an elliptic curve of order r, and P a generator of $E(\mathbb{F})$. Then the elliptic curve scalar multiplication with base P is defined as follows (where $[0]P = \mathcal{O}$ and [m]P = P + P + ... + P is the *m*-fold sum of *P* with itself):

$$[\cdot]P:\mathbb{Z}_r\to E(\mathbb{F}); m\mapsto [m]P$$

therefore, elliptic curve scalar multiplication is an instantiation of the general exponential map using additive instead of multiplicative notation. This map is a homomorph of groups, which means that $[n+m]P = [n]P \oplus [m]P$.

As with all finite, cyclic groups, the inverse of the exponential map exists and is usually called the **elliptic curve discrete logarithm map**. However, elliptic curves are believed to be XXX-groups, which means that we don't know of any efficient way to actually compute this

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Scalar multiplication and its inverse, the elliptic curve discrete logarithm, define the elliptic curve discrete logarithm problem, which consists of finding solutions $m \in \mathbb{Z}_r$ such that the following equation holds:

$$P = [m]Q \tag{5.3}$$

Any solution m is usually called a **discrete logarithm relation** between P and Q. If Q is a generator of the curve, then there is a discrete logarithm relation between Q and any other point, since Q generates the group by repeatedly adding Q to itself. Therefore, we know that some discrete logarithm relation exists for generator Q and point P. However, since elliptic curves are believed to be XXX-groups, finding actual relations m is computationally hard, with runtimes being approximately the size of the order of the group. In practice, we often need the assumption that a discrete logarithm relation exists, while the relation itself is not known.

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One useful property of the exponential map in regard to the examples in this book is that it can be used to greatly simplify pen-and-paper computations. As we have seen in example XXX, computing the elliptic curve addition law takes quite a bit of effort when done without a computer. However, when g is a generator of a small pen-and-paper elliptic curve group of order r, we can use the exponential map to write the group using cofactor clearing, which implies that $[r]g = \mathscr{O}$:

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cofactor clearing

$$\mathbb{G} = \{ [1]g \to [2]g \to [3]g \to \dots \to [r-1]g \to \mathcal{O} \}$$
 (5.4)

"Logarithmic ordering" like this greatly simplifies complicated elliptic curve addition to the much simpler case of modular r addition. In order to add two curve points P and Q, we only have to look up their discrete log relations with the generator, say P = [n]g and Q = [m]g, and compute the sum as $P \oplus Q = [n+m]g$. This is, of course, only possible for small groups where we can keep a clear overview, such as XXX.

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In the following example, we will look at some implications of the fact that elliptic curves are finite cyclic groups. We will apply the fundamental theorem of finite cyclic groups and look how it reflects on the curves in consideration.

Example 73. Consider the elliptic curve group $E_1(\mathbb{F}_5)$ from example 65. Since it is a finite cyclic group of order 9, and the prime factorization of 9 is $3 \cdot 3$, we can use the fundamental

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theorem of finite cyclic groups to reason about all its subgroups. In fact, since the only prime factor of 9 is 3, we know that $E_1(\mathbb{F}_5)$ has the following subgroups:

- $\mathbb{G}_1 = E_1(\mathbb{F}_5)$ is a subgroup of order 9. By definition, any group is a subgroup of itself.
- $\mathbb{G}_2 = \{(2,1),(2,4),\mathcal{O}\}$ is a subgroup of order 3. This is the subgroup associated to the prime factor 3.
- $\mathbb{G}_3 = \{ \mathcal{O} \}$ is a subgroup of order 1. This is the trivial subgroup.

Moreover, since $E_1(\mathbb{F}_5)$ and all its subgroups are cyclic, we know from page 43 that they must have generators. For example, the curve point (2,1) is a generator of the order 3 subgroup \mathbb{G}_2 , since every element of \mathbb{G}_2 can be generated by repeatedly adding (2,1) to itself:

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$$[1](2,1) = (2,1)$$
$$[2](2,1) = (2,4)$$
$$[3](2,1) = \emptyset$$

Since (2,1) is a generator, we know from XXX that it gives rise to an exponential map from the finite field \mathbb{F}_3 onto \mathbb{G}_2 defined by scalar multiplication:

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$$[\cdot](2,1): \mathbb{F}_3 \to \mathbb{G}_2: x \mapsto [x](2,1)$$

To give an example of a generator that generates the entire group $E_1(\mathbb{F}_5)$, consider the point (0,1). Applying the tangent rule repeatedly, we compute as follows:

Again, since (2,1) is a generator, we know from XXX that it gives rise to an exponential map. However, since the group order is not a prime number, the exponential map does not map from any field, but from the residue class ring \mathbb{Z}_9 only:

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$$[\cdot](0,1): \mathbb{Z}_9 \to \mathbb{G}_1 : x \mapsto [x](0,1)$$

Using the generator (0,1) and its associated exponential map, we can write $E(\mathbb{F}_1)$ i logarithmic order with respect to (0,1) as explained in equation 5.4. We get the following:

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$$E_1(\mathbb{F}_5) = \{(0,1) \to (4,2) \to (2,1) \to (3,4) \to (3,1) \to (2,4) \to (4,3) \to (0,4) \to \mathscr{O}\}$$

This indicates that the first element is a generator, and the n-th element is the scalar product of n and the generator. To see how logarithmic orders like this simplify the computations in small elliptic curve groups, consider example 70 again. In that example, we use the chord-and-tangent rule to compute $(0,1) \oplus (4,2)$. Now, in the logarithmic order of $E_1(\mathbb{F})$, we can compute that sum much easier, since we can directly see that (0,1) = [1](0,1) and (4,2) = [2](0,1). We can then deduce $(0,1) \oplus (4,2) = (2,1)$ immediately, since $[1](0,1) \oplus [2](0,1) = [3](0,1) = (2,1)$.

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To give another example, we can immediately see that $(3,4) \oplus (4,3) = (4,2)$, without doing any expensive elliptic curve addition, since we know (3,4) = [4](0,1) as well as (4,3) = (4,3)

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[7](0,1) from the logarithmic representation of $E_1(\mathbb{F}_5)$. Since 4+7=2 in \mathbb{Z}_9 , the result must be [2](0,1) = (4,2).

Finally we can use $E_1(\mathbb{F}_5)$ as an example to understand the concept of cofactor clearing from 5.4. Since the order of $E_1(\mathbb{F}_5)$ is 9, we only have a single factor, which happen to be the cofactor check as well. Cofactor clearing then implies that we can map any element from $E_1(\mathbb{F}_5)$ onto its prime factor group \mathbb{G}_2 by scalar multiplication with 3. For example, taking the element (3,4), which is not in \mathbb{G}_2 , and multiplying it with 3, we get [3](3,4) = (2,1), which is an element of \mathbb{G}_2 as expected.

reference

In the following example, we will look at the subgroups of our pen-jubjub curve, define generators, and compute the logarithmic order for pen-and-paper computations. Then we take another look at the principle of cofactor clearing.

Example 74. Consider the pen-jubjub curve PJJ_13 from example 66 again. Since the order of PJJ_13 is 20, and the prime factorization of 20 is $2^2 \cdot 5$, we know that the PJJ_13 contains a "large" prime-order subgroup of size 5 and a small prime oder subgroup of size 2.

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To compute those groups, we can apply the technique of cofactor clearing in a try-and-repeat loop. We start the loop by arbitrarily choosing an element $P \in PJJ$ 13, then multiplying that element with the cofactor of the group that we want to compute. If the result is \mathcal{O} , we try a different element and repeat the process until the result is different from the point at infinity \mathcal{O} .

To compute a generator for the small prime-order subgroup (PJJ_13)₂, first observe that the cofactor is 10, since $20 = 2 \cdot 10$. We then arbitrarily choose the curve point $(5,11) \in PJJ_13$ and compute $[10](5,11) = \emptyset$. Since the result is the point at infinity, we have to try another curve point, say (9,4). We get [10](9,4) = (4,0) and we can deduce that (4,0) is a generator of $(PJJ_13)_2$. Logarithmic order then gives $(PJJ_13)_2 = \{(4,0) \to \emptyset\}$ as expected, since we know from example 71 that (4,0) is self-inverse, with $(4,0) \oplus (4,0) = \emptyset$. We double check the computations using Sage:

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```
sage: F13 = GF(13)
                                                                                    304
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    sage: PJJ = EllipticCurve(F13,[8,8])
                                                                                    305
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    sage: P = PJJ(5,11)
                                                                                    306
2912
    sage: INF = PJJ(0)
                                                                                    307
    sage: 10*P == INF
                                                                                    308
2914
    True
                                                                                    309
2915
    sage: Q = PJJ(9,4)
                                                                                    310
2916
    sage: R = PJJ(4,0)
                                                                                    311
2917
    sage: 10*Q == R
                                                                                    312
2918
    True
                                                                                    313
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```

We can apply the same reasoning to the "large" prime-order subgroup (PJJ_13)₅, which contains 5 elements. To compute a generator for this group, first observe that the associated cofactor is 4, since 20 = 5.4. We choose the curve point $(9,4) \in PJJ$ 13 again, and compute [4](9,4) = (7,11). We can deduce that (7,11) is a generator of $(PJJ_13)_5$. Using the generator Explain (7,11), we compute the exponential map $[\cdot](7,11): \mathbb{F}_5 \to PJJ_13$ and get the following:

how

```
[0](7,11) = \mathcal{O}
[1](7,11) = (7,11)
[2](7,11) = (8,5)
[3](7,11) = (8,8)
[4](7,11) = (7,2)
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We can use this computation to write the large-order prime group (PJJ_13)₅ of the pen-jubjub 2920 curve in logarithmic order, which we will use quite frequently in what follows. We get the following:

$$(PJJ_13)_5 = \{(7,11) \to (8,5) \to (8,8) \to (7,2) \to \mathcal{O}\}\$$
 (5.5)

From this, we can immediately see, for example that $(8,8) \oplus (7,2) = (8,5)$, since 3+4=2 in 2923 \mathbb{F}_5 . 2924

From the previous two examples, the reader might get the impression that elliptic curve computation can be largely replaced by modular arithmetics. This however, is not true in general, but only an artifact of small groups, where it is possible to write the entire group in a logarithmic order. The following example gives some understanding of why this is not possible in cryptographically secure groups.

Example 75. SEKTP BICOIN. DISCRETE LOG HARDNESS PROHIBITS ADDITION IN THE FIELD...

write example

Projective short Weierstraß form As we have seen in the previous section, describing elliptic curves as pairs of points that satisfy a certain equation is relatively straight-forward. However, in order to define a group structure on the set of points, we had to add a special point at infinity to act as the neutral element.

Recalling from the definition of projective planes (section 4.4), we know that points at infinity are handled as ordinary points in projective geometry. Therefore, it makes sense to look at the definition of a short Weierstraß curve in projective geometry.

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To see what a short Weierstraß curve in projective coordinates is, let $\mathbb F$ be a finite field of order q and characteristic > 3, let $a, b \in \mathbb{F}$ be two field elements such that $4a^3 + 27b^2 \mod q \neq 0$ and let $\mathbb{F}P^2$ be the projective plane over \mathbb{F} . Then a short Weierstraß elliptic curve over \mathbb{F} in its projective representation is the set of all points $[X:Y:Z] \in \mathbb{F}P^2$ from the projective plane that satisfy the **homogenous** cubic equation $Y^2 \cdot Z = X^3 + a \cdot X \cdot Z^2 + b \cdot Z^3$:

$$E(\mathbb{F}P^2) = \{ [X : Y : Z] \in \mathbb{F}P^2 \mid Y^2 \cdot Z = X^3 + a \cdot X \cdot Z^2 + b \cdot Z^3 \}$$
 (5.6)

To understand how the point at infinity is unified in this definition, recall from XXX that, in projective geometry, points at infinity are given by homogeneous coordinates [X:Y:0]. Inserting representatives $(x_1, y_1, 0) \in [X : Y : 0]$ from those classes into the defining homogenous cubic equations gives the following:

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$$y_1^2 \cdot 0 = x_1^3 + a \cdot x_1 \cdot 0^2 + b \cdot 0^3 \qquad \Leftrightarrow 0 = x_1^3$$

This shows that the only point at infinity, that is also a point on a projective short Weierstraß curve is the class $[0, 1, 0] = \{(0, y, 0) \mid y \in \mathbb{F}\}.$

This point is the projective representation of \mathcal{O} . The projective representation of a short Weierstraß curve, therefore, has the advantage that it does not need a special symbol to represent the point at infinity \mathcal{O} from the affine definition.

Example 76. To get an intuition of how an elliptic curve in projective geometry looks, consider curve $E_1(\mathbb{F}_5)$ from example (65). We know that, in its affine representation, the set of rational points is given as follows:

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$$E_1(\mathbb{F}_5) = \{ \mathscr{O}, (0,1), (2,1), (3,1), (4,2), (4,3), (0,4), (2,4), (3,4) \}$$
 (5.7)

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This is defined as the set of all pairs $(x,y) \in \mathbb{F}_5 \times \mathbb{F}_5$ such that the affine short Weierstraß equation $y^2 = x^3 + ax + b$ with a = 1 and b = 1 is satisfied.

To find the projective representation of a short Weierstraß curve with the same parameters a=1 and b=1, we have to compute the set of projective points [X:Y:Z] from the projective plane $\mathbb{F}_5\mathsf{P}^2$ that satisfy the following homogenous cubic equation for any representative $(x_1,y_1,z_1)\in [X:Y:Z]$:

$$y_1^2 z_1 = x_1^3 + 1 \cdot x_1 z_1^2 + 1 \cdot z_1^3 \tag{5.8}$$

We know from XXX that the projective plane \mathbb{F}_5P^2 contains $5^2 + 5 + 1 = 31$ elements, so we can take the effort and insert all elements into equation 5.8 and see if both sides match.

For example, consider the projective point [0:4:1]. We know from XXX that this point in the projective plane represents the following line in the three-dimensional space \mathbb{F}^3 :

$$[0:4:1] = \{(0,0,0), (0,4,1), (0,3,2), (0,2,3), (0,1,4)\}$$

To check whether or not [0:4:1] satisfies 5.8, we can insert any representative, in other words, any element from XXX. Each element satisfies the equation if and only if all other elements satisfy the equation. So, we insert (0,4,1) and get the following result:

$$1^2 \cdot 1 = 0^3 + 1 \cdot 0 \cdot 1^2 + 1 \cdot 1^3$$

This tells us that the affine point [0:4:1] is indeed a solution to the equation 5.8, but we could just as well have inserted any other representative. For example, inserting (0,3,2) also satisfies 5.8:

$$3^2 \cdot 2 = 0^3 + 1 \cdot 0 \cdot 2^2 + 1 \cdot 2^3$$

To find the projective representation of E_1 , we first observe that the projective line at infinity [1:0:0] is not a curve point on any projective short Weierstraß curve, since it cannot satisfy XXX for any parameter a and b. Therefore, we can exclude it from our consideration.

Moreover, a point at infinity [X : Y : 0] can only satisfy equation XXX for any a and b, if X = 0, which implies that the only point at infinity relevant for short Weierstraß elliptic curves is [0:1:0], since [0:k:0] = [0:1:0] for all k from the finite field. Therefore, we can exclude all points at infinity except the point [0:1:0].

All points that remain are the affine points [X : Y : 1]. Inserting all of them into XXX, we get the set of all projective curve points as follows:

$$E_1(\mathbb{F}_5 \mathbf{P}^2) = \{ [0:1:0], [0:1:1], [2:1:1], [3:1:1], \\ [4:2:1], [4:3:1], [0:4:1], [2:4:1], [3:4:1] \}$$

If we compare this with the affine representation, we see that there is a 1:1 correspondence between the points in the affine representation in 5.7 and the affine points in projective geometry, and that the point [0:1:0] represents the additional point \mathcal{O} in the projective representation.

Exercise 37. Compute the projective representation of the pen-jubjub curve and the logarithmic order of its large prime-order subgroup with respect to the generator (7, 11).

Projective Group law As we have seen on page 70, one of the key properties of an elliptic curve is that it comes with a definition of a group law on the set of its rational points, described geometrically by the chord-and-tangent rule (definition 5.1.1.1). This rule was kind of intuitive,

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with the exception of the distinguished point at infinity, which appeared whenever the chord or the tangent did not have a third intersection point with the curve.

One of the key features of projective coordinates is that, in projective space, it is guaranteed that any chord will always intersect the curve in three points, and any tangent will intersect it in two points including the tangent point. So, the geometric picture simplifies, as we don't need to consider external symbols and associated cases.

Again, it can be shown that the points of an elliptic curve in projective space form a commutative group with respect to the tangent-and-chord rule such that the projective point [0:1:0] is the neutral element, and the additive inverse of a point [X:Y:Z] is given by [X:-Y:Z]. The addition law is usually described by the following algorithm, minimizing the number of necessary additions and multiplications in the base field.

Exercise 38. Compare the affine addition law for short Weierstraß curves with the projective addition rule. Which branch in the projective rule corresponds to which case in the affine law?

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Coordinate Transformations As we have seen in example XXX, there was a close relation between the affine and the projective representation of a short Weierstraß curve. This was not a coincidence. In fact, from a mathematical point of view, projective and affine short Weierstraß curves describe the same thing, as there is a one-to-one correspondence (an isomorphism) between both representations for any arbitrary parameters a and b.

To specify the isomorphism, let $E(\mathbb{F})$ and $E(\mathbb{F}P^2)$ be an affine and a projective short Weierstraß curve defined for the same parameters a and b. Then the map in 5.9 maps points from the affine representation to points from the projective representation of a short Weierstraß curve. In other words, if the pair of points (x,y) satisfies the affine equation $y^2 = x^3 + ax + b$, then all homogeneous coordinates $(x_1,y_1,z_1) \in [x:y:1]$ satisfy the projective equation $y_1^2 \cdot z_1 = x_1^3 + ay_1 \cdot z_1^2 + b \cdot z_1^3$.

$$\Phi: E(\mathbb{F}) \to E(\mathbb{F}P^2) : \begin{array}{ccc} (x,y) & \mapsto & [x:y:1] \\ \mathscr{O} & \mapsto & [0:1:0] \end{array}$$
 (5.9)

The inverse is given by the following map:

$$\Phi^{-1}: E(\mathbb{FP}^2) \to E(\mathbb{F}) : [X:Y:Z] \mapsto \begin{cases} (\frac{X}{Z}, \frac{Y}{Z}) & \text{if } Z \neq 0\\ \emptyset & \text{if } Z = 0 \end{cases}$$
 (5.10)

Note that the only projective point [X:Y:Z] with $Z \neq 0$ that satisfies XXX is given by the class [0:1:0].

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One key feature of Φ and its inverse is that it respects the group structure, which means that $\Phi((x_1,y_1)\oplus(x_2,y_2))$ is equal to $\Phi(x_1,y_1)\oplus\Phi(x_2,y_2)$. The same holds true for the inverse map Φ^{-1}

Maps with these properties are called **group isomorphisms**, and, from a mathematical point of view, the existence of Φ implies that these two definitions are equivalent, and implementations can choose freely between these representations.

5.1.2 Montgomery Curves

History and use of them (optimized scalar multiplication)

write up this part

Algorithm 6 Projective Weierstraß Addition Law

```
Require: [X_1:Y_1:Z_1], [X_2:Y_2:Z_2] \in E(\mathbb{FP}^2)
   procedure ADD-RULE([X_1 : Y_1 : Z_1], [X_2 : Y_2 : Z_2])
         if [X_1:Y_1:Z_1] == [0:1:0] then
              [X_3:Y_3:Z_3] \leftarrow [X_2:Y_2:Z_2]
         else if [X_2:Y_2:Z_2] == [0:1:0] then
              [X_3:Y_3:Z_3] \leftarrow [X_1:Y_1:Z_1]
         else
              U_1 \leftarrow Y_2 \cdot Z_1
              U_2 \leftarrow Y_1 \cdot Z_2
              V_1 \leftarrow X_2 \cdot Z_1
              V_2 \leftarrow X_1 \cdot Z_2
              if V_1 == V_2 then
                   if U_1 \neq U_2 then [X_3 : Y_3 : Z_3] \leftarrow [0 : 1 : 0]
                         if Y_1 == 0 then [X_3 : Y_3 : Z_3] \leftarrow [0 : 1 : 0]
                         else
                              W \leftarrow a \cdot Z_1^2 + 3 \cdot X_1^2
                              S \leftarrow Y_1 \cdot Z_1
                              B \leftarrow X_1 \cdot Y_1 \cdot S
                              H \leftarrow W^2 - 8 \cdot B
                              X' \leftarrow 2 \cdot H \cdot S
                              Y' \leftarrow W \cdot (4 \cdot B - H) - 8 \cdot Y_1^2 \cdot S^2
                              Z' \leftarrow 8 \cdot S^3
                              [X_3:Y_3:Z_3] \leftarrow [X':Y':Z']
                         end if
                   end if
              else
                    U = U_1 - U_2
                    V = V_1 - V_2
                   W = Z_1 \cdot Z_2
                   A = U^{2} \cdot W - V^{3} - 2 \cdot V^{2} \cdot V_{2}
                   X' = V \cdot A
                   Y' = U \cdot (V^2 \cdot V_2 - A) - V^3 \cdot U_2
                   Z' = V^3 \cdot W
                   [X_3:Y_3:Z_3] \leftarrow [X':Y':Z']
              end if
         end if
         return [X_3 : Y_3 : Z_3]
   end procedure
Ensure: [X_3:Y_3:Z_3] == [X_1:Y_1:Z_1] \oplus [X_2:Y_2:Z_2]
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Affine Montgomery Form To see what a Montgomery curve in affine coordinates is, let \mathbb{F} be a finite field of characteristic > 2, and let $A, B \in \mathbb{F}$ be two field elements such that $B \neq 0$ and $A^2 \neq 4$. A Montgomery elliptic curve $M(\mathbb{F})$ over \mathbb{F} in its affine representation is the set of all pairs of field elements $(x,y) \in \mathbb{F} \times \mathbb{F}$ that satisfy the Montgomery cubic equation $B \cdot y^2 = x^3 + A \cdot x^2 + x$, together with a distinguished symbol \mathcal{O} , called the **point at infinity**.

$$M(\mathbb{F}) = \{(x, y) \in \mathbb{F} \times \mathbb{F} \mid B \cdot y^2 = x^3 + A \cdot x^2 + x\} \bigcup \{\emptyset\}$$
 (5.11)

Despite the fact that Montgomery curves look different from short Weierstraß curves, they are just a special way to describe certain short Weierstraß curves. In fact, every curve in affine Montgomery form can be transformed into an elliptic curve in Weierstraß form. To see that, assume that a curve is given in Montgomery form $By^2 = x^3 + Ax^2 + x$. The associated Weierstraß form is then as follows:

is the label in LATEX correct here?

$$y^2 = x^3 + \frac{3 - A^2}{3B^2} \cdot x + \frac{2A^3 - 9A}{27B^3}$$
 (5.12)

On the other hand, an elliptic curve $E(\mathbb{F})$ over base field \mathbb{F} in Weierstraß form $y^2 = x^3 + ax + b$ can be converted to Montgomery form if and only if the following conditions hold:

Definition 5.1.2.1. Requirements for Montgomery curves

- The number of points on E(F) is divisible by 4
- The polynomial $z^3 + az + b \in \mathbb{F}[z]$ has at least one root $z_0 \in \mathbb{F}$
- $3z_0^2 + a$ is a quadratic residue in \mathbb{F} .

When these conditions are satisfied, then for $s = (\sqrt{3z_0^2 + a})^{-1}$, the equivalent Montgomery curve is defined by the following equation:

$$sy^2 = x^3 + (3z_0s)x^2 + x (5.13)$$

In the following example we will look at our pen-jubjub curve again, and show that it is actually a Montgomery curve.

Example 77. Consider the prime field \mathbb{F}_{13} and the pen-jubjub curve PJJ_13 from example 66. To see that it is a Montgomery curve, we have to check the requirements from 5.1.2.1:

Since the order of *PJJ_13* is 20, which is divisible by 4, the first requirement is met.

Next, since a=8 and b=8, we have to check if the polynomial $P(z)=z^3+8z+8$ has a root in \mathbb{F}_{13} . We simply evaluate P at all numbers $z \in \mathbb{F}_{13}$, and find that P(4)=0, so a root is given by $z_0=4$.

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In the last step, we have to check that $3 \cdot z_0^2 + a$ has a root in \mathbb{F}_{13} . We compute as follows:

$$3z_0^2 + a = 3 \cdot 4^2 + 8$$

$$= 3 \cdot 3 + 8$$

$$= 9 + 8$$

$$= 4$$

To see if 4 is a quadratic residue, we can use Euler's criterion (4.28) to compute the Legendre symbol of 4. We get the following:

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$$\left(\frac{4}{13}\right) = 4^{\frac{13-1}{2}} = 4^6 = 1$$

This means that 4 does have a root in \mathbb{F}_{13} . In fact, computing a root of 4 in \mathbb{F}_{13} is easy, since the integer root 2 of 4 is also one of its roots in \mathbb{F}_{13} . The other root is given by 13 - 4 = 9.

Since all requirements are meet, we have now shown that *PJJ_13* is indeed a Montgomery curve, and we can use 5.13 to compute its associated Montgomery form. We compute as follows:

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$$s = \left(\sqrt{3 \cdot z_0^2 + 8}\right)^{-1}$$

$$= 2^{-1}$$

$$= 2^{13-2}$$

$$= 7$$
Fermat's little theorem
$$= 2048 \mod 13 = 7$$

The defining equation for the Montgomery form of our pen-jubjub curve is then given by the following equation:

$$sy^{2} = x^{3} + (3z_{0}s)x^{2} + x \qquad \Rightarrow$$

$$7 \cdot y^{2} = x^{3} + (3 \cdot 4 \cdot 7)x^{2} + x \qquad \Leftrightarrow$$

$$7 \cdot y^{2} = x^{3} + 6x^{2} + x$$

So, we get the defining parameters as B = 7 and A = 6, and we can write the pen-jubjub curve in its affine Montgomery representation as follows:

$$PJJ_{13} = \{(x,y) \in \mathbb{F}_{13} \times \mathbb{F}_{13} \mid 7 \cdot y^2 = x^3 + 6x^2 + x\} \bigcup \{\mathscr{O}\}$$
 (5.14)

Now that we have the abstract definition of our pen-jubjub curve in Montgomery form, we can compute the set of points by inserting all pairs $(x, y) \in \mathbb{F}_{13} \times \mathbb{F}_{13}$ similarly to how we computed the curve points in its Weierstraß representation. We get the following:

$$PJJ_{1} = \{ \mathscr{O}, (0,0), (1,4), (1,9), (2,4), (2,9), (3,5), (3,8), (4,4), (4,9), (5,1), (5,12), (7,1), (7,12), (8,1), (8,12), (9,2), (9,11), (10,3), (10,10) \}$$

```
sage: F13 = GF(13)
                                                                                314
3046
    sage: L_MPJJ = []
                                                                                315
3047
     ....: for x in F13:
                                                                                316
3048
                for y in F13:
                                                                                317
                     if F13(7)*y^2 == x^3 + F13(6)*x^2 +x:
                                                                                318
3050
     . . . . :
                          L_MPJJ.append((x,y))
                                                                                319
3051
    sage: MPJJ = Set(L_MPJJ)
                                                                                320
3052
    sage: # does not compute the point at infinity
                                                                                321
3053
```

Affine Montgomery coordinate transformation Comparing the Montgomery representation of the previous example (equation 5.14) with the Weierstraß representation of the same curve (equation 5.2), we see that there is a 1:1 correspondence between the curve points in both

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examples. This is no accident. In fact, if $M_{A,B}$ is a Montgomery curve, and $E_{a,b}$ a Weierstraß curve with $a=\frac{3-A^2}{3B^2}$ and $b=\frac{2A^2-9A}{27B^3}$ then the following function maps all points in Montgomery representation onto the points in Weierstraß representation: 3057 3059

$$\Phi: M_{A,B} \to E_{a,b} : (x,y) \mapsto \left(\frac{3x+A}{3B}, \frac{y}{B}\right)$$
 (5.15)

This map is a 1:1 correspondence (am isomorphism), and its inverse map is given by the fol-3060 lowing equation (where z_0 is a root of the polynomial $z^3 + az + b \in \mathbb{F}[z]$ and $s = (\sqrt{3z_0^2 + a})^{-1}$). 3061

$$\Phi^{-1}: E_{a,b} \to M_{A,B}: (x,y) \mapsto (s \cdot (x-z_0), s \cdot y)$$
 (5.16)

Using this map, it is therefore possible for implementations of Montgomery curves to freely 3063 transit between the Weierstraß and the Montgomery representation. However, as we saw in 3064 definition 5.1.2.1, not every Weierstraß curve is a Montgomery curve, as all criteria in 5.1.2.1 3065 have to be satisfied. This means that the map Φ^{-1} does not always exist.

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Example 78. Consider our pen-jubjub curve again. In equation 5.2 we derived its Weierstraß representation and in example 5.14, we derived its Montgomery representation.

To see how coordinate transformation Φ works in this example, let's map points from the Montgomery representation onto points from the Weierstraß representation. Inserting, for example, the point (0,0) from the Montgomery representation 5.14 into Φ gives the following:

$$\Phi(0,0) = \left(\frac{3 \cdot 0 + A}{3B}, \frac{0}{B}\right)$$
$$= \left(\frac{3 \cdot 0 + 6}{3 \cdot 7}, \frac{0}{7}\right)$$
$$= \left(\frac{6}{8}, 0\right)$$
$$= (4,0)$$

As we can see, the Montgomery point (0,0) maps to the self-inverse point (4,0) of the Weierstraß representation. On the other hand, we can use our computations of s = 7 and $z_0 = 4$ from XXX to compute the inverse map Φ^{-1} , which maps points on the Weiertraß representation to points on the Mongomery form. Inserting, for example, (4,0) we get the following:

add reference

$$\Phi^{-1}(4,0) = (s \cdot (4-z_0), s \cdot 0)$$
$$= (7 \cdot (4-4), 0)$$
$$= (0,0)$$

As expected, the inverse map maps the Weierstraß point back to where it originated in the 3069 Montgomery form. We can invoke Sage to check that our computation of Φ is correct: 3070

```
sage: # Compute PHI of Montgomery form:
                                                                              322
3071
    sage: L_PHI_MPJJ = []
                                                                              323
3072
    sage: for (x,y) in L_MPJJ: # LMJJ as defined previously
                                                                              324
3073
                v = (F13(3)*x + F13(6))/(F13(3)*F13(7))
                                                                              325
3074
                w = y/F13(7)
     . . . . :
                                                                              326
3075
                L_PHI_MPJJ.append((v,w))
3076
                                                                              327
```

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```
sage: PHI_MPJJ = Set(L_PHI_MPJJ)
                                                                               328
3077
    sage: # Computation Weierstrass form
                                                                                329
3078
    sage: C_WPJJ = EllipticCurve(F13,[8,8])
                                                                                330
3079
    sage: L_WPJJ = [P.xy() for P in C_WPJJ.points() if P.order() >
                                                                               331
3080
         11
3081
    sage: WPJJ = Set(L_WPJJ)
                                                                               332
3082
    sage: # check PHI(Montgomery) == Weierstrass
                                                                               333
3083
    sage: WPJJ == PHI_MPJJ
                                                                               334
3084
    True
                                                                               335
3085
    sage: # check the inverse map PHI^(-1)
                                                                               336
3086
    sage: L PHIINV WPJJ = []
                                                                               337
3087
    sage: for (v,w) in L_WPJJ:
                                                                               338
3088
    . . . . :
                x = F13(7)*(v-F13(4))
                                                                               339
3089
                y = F13(7) *w
                                                                               340
     . . . . :
3090
                L_PHIINV_WPJJ.append((x,y))
                                                                               341
3091
    sage: PHIINV_WPJJ = Set(L_PHIINV_WPJJ)
                                                                               342
3092
    sage: MPJJ == PHIINV WPJJ
                                                                               343
3093
    True
                                                                               344
3094
```

Montgomery group law We have seen that Montgomery curves special cases of short Weierstraß curves. As such, they have a group structure defined on the set of their points, which can also be derived from the chord-and-tangent rule. In accordance with short Weierstraß curves, it can be shown that the identity $x_1 = x_2$ implies $y_2 = \pm y_1$, meaning that the following rules are a complete description of the affine addition law.

3100 Definition 5.1.2.2. Montgomery group law

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- (Neutral element) Point at infinity \mathcal{O} is the neutral element.
- (Additive inverse) The additive inverse of \mathscr{O} is \mathscr{O} . For any other curve point $(x,y) \in M(\mathbb{F}_q) \setminus \{\mathscr{O}\}$, the additive inverse is given by (x,-y).
 - (Addition rule) For any two curve points $P, Q \in M(\mathbb{F}_q)$, addition is defined by one of the following cases:
 - 1. (Adding the neutral element) If $Q = \mathcal{O}$, then the sum is defined as P + Q = P.
 - 2. (Adding inverse elements) If P = (x, y) and Q = (x, -y), the sum is defined as $P + Q = \emptyset$.
 - 3. (Adding non-self-inverse equal points) If P = (x, y) and Q = (x, y) with $y \ne 0$, the sum 2P = (x', y') is defined as follows:

$$x' = \left(\frac{3x_1^2 + 2Ax_1 + 1}{2By_1}\right)^2 \cdot B - (x_1 + x_2) - A$$
 , $y' = \frac{3x_1^2 + 2Ax_1 + 1}{2By_1}(x_1 - x') - y_1$

4. (Adding non-inverse different points) If $P = (x_1, y_1)$ and $Q = (x_2, y_2)$ such that $x_1 \neq x_2$, the sum R = P + Q with $R = (x_3, y_3)$ is defined as follows:

$$x' = (\frac{y_2 - y_1}{x_2 - x_1})^2 B - (x_1 + x_2) - A$$
, $y' = \frac{y_2 - y_1}{x_2 - x_1} (x_1 - x') - y_1$

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5.1.3 Twisted Edwards Curves

As we have seen in 5.1.2.2 both Weierstraß and Montgomery curves have somewhat complicated addition and doubling laws, as many cases have to be distinguished. Those various cases translate to branches in computer programs.

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In the context of SNARK development, two computational models for bounded computations are used, called **circuits** and **rank-1 constraint systems**. Program branches are undesirably costly when implemented in those models. It is therefore advantageous to look for curves with an addition/doubling rule that requires no branches and as few field operations as possible.

Twisted Edwards curves are particularly useful here, as a subclass of these curves has a compact and easily implementable addition law that works for all points including the point at infinity. Implementing this law needs no branching.

Twisted Edwards Form To see what an affine **twisted Edwards curve** looks like, let \mathbb{F} be a finite field of characteristic > 2, and let $a, d \in \mathbb{F} \setminus \{0\}$ be two non-zero field elements with $a \neq d$. A **twisted Edwards elliptic curve** in its affine representation is the set of all pairs (x, y) from $\mathbb{F} \times \mathbb{F}$ that satisfy the twisted Edwards equation $a \cdot x^2 + y^2 = 1 + d \cdot x^2 y^2$, given below:

$$E(\mathbb{F}) = \{ (x, y) \in \mathbb{F} \times \mathbb{F} \mid a \cdot x^2 + y^2 = 1 + d \cdot x^2 y^2 \}$$
 (5.17)

A twisted Edwards curve is called an **Edwards curve** (**non-twisted**), if the parameter a is equal to 1, and it is called a **SNARK-friendly twisted Edwards curve** if the parameter a is a quadratic residue and the parameter d is a quadratic non-residue.

As we can see from the definition, affine twisted Edwards curves look somewhat different from Weierstraß curves, as their affine representation does not need a special symbol to represent the point at infinity. In fact, we we will see that the pair (0,1) is always a point on any twisted Edwards curve, and that it takes the role of the point at infinity.

Despite their different appearances however, twisted Edwards curves are equivalent to Montgomery curves in the sense that, for every twisted Edwards curve, there is a Montgomery curve, and a way to map the points of one curve in a 1:1 correspondence onto the other and vice versa. To see that, assume that a curve in twisted Edwards form $a \cdot x^2 + y^2 = 1 + d \cdot x^2 y^2$ is given. The associated Montgomery curve is then defined by the Montgomery equation:

$$\frac{4}{a-d}y^2 = x^3 + \frac{2(a+d)}{a-d} \cdot x^2 + x \tag{5.18}$$

On the other hand, a Montgomery curve $By^2 = x^3 + Ax^2 + x$ with $B \neq 0$ and $A^2 \neq 4$ can give rise to a twisted Edwards curve defined by the following equation:

$$\left(\frac{A+2}{B}\right)x^2 + y^2 = 1 + \left(\frac{A-2}{B}\right)x^2y^2 \tag{5.19}$$

As we have seen in equation 5.12 and the following discussion, Montgomery curves are just a special class of Weierstraß curves. Furthermore we now know that twisted Edwards curves are special Weierstraß curves too. This means that the more general way to describe elliptic curves is as Weierstraß curves.

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Example 79. Consider the pen-jubjub curve from example 66 again. We know from example 77 that it is a Montgomery curve, and, since Montgomery curves are equivalent to twisted Edwards curves, we want to write this curve in twisted Edwards form. We use equation 5.19,

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and compute the parameters a and d as follows:

$$a = \frac{A+2}{B}$$
 # insert A=6 and B=7
$$= \frac{8}{7} = 3$$
 # $7^{-1} = 2$

$$d = \frac{A-2}{B}$$

$$= \frac{4}{7} = 8$$

Thus, we get the defining parameters as a=3 and d=8. Since our goal is to use this curve later on in implementations of pen-and-paper SNARKs, let us show that tiny-jubjub is also a **SNARK-friendly** twisted Edwards curve. To see that, we have to show that a is a quadratic residue and d is a quadratic non-residue. We therefore compute the Legendre symbols of a and d using Euler's criterion. We get the following:

change
"tinyjubjub"
to "penjubjub"
throughout?

$$\left(\frac{3}{13}\right) = 3^{\frac{13-1}{2}} = 3^6 = 1$$

$$\left(\frac{8}{13}\right) = 8^{\frac{13-1}{2}}$$
$$= 8^6 = 12 = -1$$

This proves that tiny-jubjub is SNARK-friendly. We can write the tiny-jubjub curve in its affine twisted Edwards representation as follows:

$$TJJ_{13} = \{(x,y) \in \mathbb{F}_{13} \times \mathbb{F}_{13} \mid 3 \cdot x^2 + y^2 = 1 + 8 \cdot x^2 \cdot y^2\}$$
 (5.20)

Now that we have the abstract definition of our pen-jubjub curve in twisted Edwards form, we can compute the set of points by inserting all pairs $(x,y) \in \mathbb{F}_{13} \times \mathbb{F}_{13}$, similarly to how we computed the curve points in its Weierstraß or Edwards representation. We get the following:

$$PJJ_{13} = \{(0,1), (0,12), (1,2), (1,11), (2,6), (2,7), (3,0), (5,5), (5,8), (6,4), (6,9), (7,4), (7,9), (8,5), (8,8), (10,0), (11,6), (11,7), (12,2), (12,11)\}$$
(5.21)

```
sage: F13 = GF(13)
                                                                                          345
3148
     sage: L_EPJJ = []
                                                                                          346
     ....: for x in F13:
3150
                                                                                          347
                  for y in F13:
                                                                                          348
3151
                        if F13(3) \times x^2 + y^2 == 1 + F13(8) \times x^2 \times y^2:
                                                                                          349
     . . . . :
3152
                             L_EPJJ.append((x,y))
                                                                                          350
3153
     sage: EPJJ = Set(L_EPJJ)
                                                                                          351
3154
```

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Twisted Edwards group law As we have seen, twisted Edwards curves are equivalent to Montgomery curves, and, as such, also have a group law. However, in contrast to Montgomery and Weierstraß curves, the group law of SNARK-friendly twisted Edwards curves can be described by a single computation that works in all cases, no matter if we add the neutral element, the inverse, or if we have to double a point. To see what the group law looks like, first observe that the point (0,1) is a solution to $a \cdot x^2 + y^2 = 1 + d \cdot x^2 \cdot y^2$ for any curve. The sum of any two points (x_1, y_1) , (x_2, y_2) on an Edwards curve $E(\mathbb{F})$ is then given by the following equation:

$$(x_1, y_1) \oplus (x_2, y_2) = \left(\frac{x_1 y_2 + y_1 x_2}{1 + dx_1 x_2 y_1 y_2}, \frac{y_1 y_2 - ax_1 x_2}{1 - dx_1 x_2 y_1 y_2}\right)$$
(5.22)

and it can be shown that the point (0,1) serves as the neutral element and the inverse of a point (x_1, y_1) is given by $(-x_1, y_1)$.

Example 80. Lets look at the tiny-jubjub curve in Edwards form from example 5.20 again. As check we have seen, this curve is given by

reference

$$PJJ_{13} = \{(0,1), (0,12), (1,2), (1,11), (2,6), (2,7), (3,0), (5,5), (5,8), (6,4), (6,9), (7,4), (7,9), (8,5), (8,8), (10,0), (11,6), (11,7), (12,2), (12,11)\}$$

To get an understanding of the twisted Edwards addition law, let's first add the neutral element (0,1) to itself. We apply the group law 5.22 and get the following:

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$$(0,1) \oplus (0,1) = \left(\frac{0 \cdot 1 + 1 \cdot 0}{1 + 8 \cdot 0 \cdot 0 \cdot 1 \cdot 1}, \frac{1 \cdot 1 - 3 \cdot 0 \cdot 0}{1 - 8 \cdot 0 \cdot 0 \cdot 1 \cdot 1} \right)$$

$$= (0,1)$$

So, as expected, adding the neutral element added to itself gives the neutral element again. Now let's add the neutral element to some other curve point. We get the following:

$$(0,1) \oplus (8,5) = \left(\frac{0 \cdot 5 + 1 \cdot 8}{1 + 8 \cdot 0 \cdot 8 \cdot 1 \cdot 5}, \frac{1 \cdot 5 - 3 \cdot 0 \cdot 8}{1 - 8 \cdot 0 \cdot 8 \cdot 1 \cdot 5}\right)$$
$$= (8,5)$$

Again, as expected, adding the neutral element to any element will result in that element again. Given any curve point (x, y), we know that its inverse is given by (-x, y). To see how the addition of a point to its inverse works, we compute as follows:

$$(5,5) \oplus (8,5) = \left(\frac{5 \cdot 5 + 5 \cdot 8}{1 + 8 \cdot 5 \cdot 8 \cdot 5 \cdot 5}, \frac{5 \cdot 5 - 3 \cdot 5 \cdot 8}{1 - 8 \cdot 5 \cdot 8 \cdot 5 \cdot 5}\right)$$

$$= \left(\frac{12 + 1}{1 + 5}, \frac{12 - 3}{1 - 5}\right)$$

$$= \left(\frac{0}{6}, \frac{12 + 10}{1 + 8}\right)$$

$$= \left(0, \frac{9}{9}\right)$$

$$= (0,1)$$

Adding a curve point to its inverse gives the neutral element, as expected. As we have seen from these examples, the twisted Edwards addition law handles edge cases particularly well and in a unified way.

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5.2 Elliptic Curve Pairings

As we have seen in equation 4.1, some groups come with the notation of a so-called pairing check map, which is a non-degenerate bilinear map from two groups into another group.

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In this section, we discuss pairings on elliptic curves, which form the basis of several zk-SNARKs and other zero-knowledge proof schemes. The SNARKs derived from pairings have the advantage of constant proof sizes, which is crucial to blockchains.

We start out by defining elliptic curve pairings and discussing a simple application which bears some resemblance to more advanced SNARKs. We then introduce the pairings arising from elliptic curves and describe Miller's algorithm, which makes these pairings practical rather than just theoretically interesting.

Elliptic curves have a few structures, like the Weil or the Tate map that qualifies as pairing.

either expand on this or delete it

Embedding Degrees As we will see in what follows, every elliptic curves gives rise to a pairing map. However, we will also see in example XXX that not every such pairing can be efficiently computed. In order to distinguish curves with efficiently computable pairings from the rest, we need to start with an introduction to the so-called **embedding degree** of a curve.

add reference

Definition 5.2.0.1. Embedding degree

Let \mathbb{F} be a finite field, let $E(\mathbb{F})$ be an elliptic curve over \mathbb{F} , and let n be a prime number that 3183 divides the order of $E(\mathbb{F})$. The embedding degree of $E(\mathbb{F})$ with respect to n is then the smallest 3184 integer k such that n divides $q^k - 1$. 3185

Fermat's little theorem (page 21 ff.) implies that every curve has at least **some** embedding degree k, since at least k = n - 1 is always a solution to the congruency $q^k \equiv 1 \pmod{n}$. This implies that the remainder of the integer division of $q^k - 1$ by n is 0.

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Example 81. To get a better intuition of the embedding degree, let's consider the elliptic curve $E_1(\mathbb{F}_5)$ from example 65. We know from 65 that the order of $E_1(\mathbb{F}_5)$ is 9, and, since the only prime factor of 9 is 3, we compute the embedding degree of $E_1(\mathbb{F}_5)$ with respect to 3.

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To find the embedding degree, we have to find the smallest integer k such that 3 divides $q^k - 1 = 5^k - 1$. We try and increment until we find a proper k.

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$$k = 1: 5^{1} - 1 = 4$$
 not divisible by 3
 $k = 2: 5^{2} - 1 = 24$ divisible by 3

Now we know that the embedding degree of $E_1(\mathbb{F}_5)$ is 2 relative to the prime factor 3.

Example 82. Let us consider the tiny jubjub curve TJJ_13 from example 66. We know from 66 that the order of TJJ 13 is 20, and that the order therefore has two prime factors. A "large" prime factor 5 and a small prime factor 2.

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We start by computing the embedding degree of TJJ_13 with respect to the large prime factor 5. To find that embedding degree, we have to find the smallest integer k such that 5 divides $q^k - 1 = 13^k - 1$. We try and increment until we find a proper k.

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$$k = 1$$
: $13^{1} - 1 = 12$ not divisible by 5
 $k = 2$: $13^{2} - 1 = 168$ not divisible by 5
 $k = 3$: $13^{3} - 1 = 2196$ not divisible by 5
 $k = 4$: $13^{4} - 1 = 28560$ divisible by 5

Now we know that the embedding degree of *TJJ_13* is 4 relative to the prime factor 5.

In real-world applications, like on pairing-friendly elliptic curves such as BLS_12-381, usually only the embedding degree of the large prime factor is relevant, which in the case of our tiny-jubjub curve is represented by 5. It should be noted, however that every prime factor of a curve's order has its own notation of embedding degree despite the fact that this is mostly irrelevant in applications.

To find the embedding degree of the small prime factor 2, we have to find the smallest integer k such that 2 divides $q^k - 1 = 13^k - 1$. We try and increment until we find a proper k.

$$k = 1$$
: $13^1 - 1 = 12$ divisible by 2

Now we know that the embedding degree of *TJJ_13* is 1 relative to the prime factor 2.

As we have seen, different prime factors can have different embedding degrees in general.

```
sage: p = 13
                                                                                      352
3206
    sage: # large prime factor
                                                                                      353
3207
     sage: n = 5
                                                                                      354
3208
     sage: for k in range(1,5): # Fermat's little theorem
                                                                                      355
3209
                 if (p^k-1) n == 0:
                                                                                      356
3210
                      break
                                                                                      357
     . . . . :
3211
    sage: k
                                                                                      358
3212
     4
                                                                                      359
3213
    sage: # small prime factor
                                                                                      360
3214
     sage: n = 2
                                                                                      361
3215
     sage: for k in range(1,2): # Fermat's little theorem
                                                                                      362
3216
                 if (p^k-1)%n == 0:
     . . . . :
                                                                                      363
3217
     . . . . :
                      break
                                                                                      364
3218
    sage: k
                                                                                      365
3219
                                                                                      366
3220
```

Example 83. To give an example of a cryptographically secure real-world elliptic curve that does not have a small embedding degree, let's look at curve Secp256k1 again. We know from 67 that the order of this curve is a prime number, so we only have a single embedding degree.

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To test potential embedding degrees k, say, in the range 1...1000, we can invoke Sage and compute as follows:

```
sage: p = 1157920892373161954235709850086879078532699846656405
3226
       64039457584007908834671663
3227
    sage: n = 1157920892373161954235709850086879078528375642790749
                                                                              368
3228
       04382605163141518161494337
3229
    sage: for k in range(1,1000):
                                                                              369
3230
                if (p^k-1) n == 0:
                                                                              370
3231
                    break
                                                                              371
    . . . . :
3232
    sage: k
                                                                              372
3233
    999
                                                                              373
3234
```

We see that Secp256k1 has at least no embedding degree k < 1000, which renders Secp256k1 a curve that has no small embedding degree. This property will be of importance later on.

Elliptic Curves over extension fields Suppose that p is a prime number, and \mathbb{F}_p its associated prime field. We know from equation 4.29 that the fields \mathbb{F}_{p^m} are extensions of \mathbb{F}_p in the sense

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that \mathbb{F}_p is a subfield of \mathbb{F}_{p^m} . This implies that we can extend the affine plane that an elliptic curve is defined on by changing the base field to any extension field. To be more precise, let $E(\mathbb{F}) = \{(x,y) \in \mathbb{F} \times \mathbb{F} \mid y^2 = x^3 + a \cdot x + b\}$ be an affine short Weierstraß curve, with parameters a and b taken from \mathbb{F} . If \mathbb{F}' is an extension field of \mathbb{F} , then we extend the domain of the curve by defining $E(\mathbb{F}')$ as follows:

$$E(\mathbb{F}') = \{ (x, y) \in \mathbb{F}' \times \mathbb{F}' \mid y^2 = x^3 + a \cdot x + b \}$$
 (5.23)

While we did not change the defining parameters, we consider curve points from the affine plane over the extension field now. Since $\mathbb{F} \subset \mathbb{F}'$, it can be shown that the original elliptic curve $E(\mathbb{F})$ is a sub-curve of the extension curve $E(\mathbb{F}')$.

Example 84. Consider the prime field \mathbb{F}_5 from example 59 and the elliptic curve $E_1(\mathbb{F}_5)$ from example 65. Since we know from XXX that \mathbb{F}_{5^2} is an extension field of \mathbb{F}_5 , we can extend the definition of $E_1(\mathbb{F}_5)$ to define a curve over \mathbb{F}_{5^2} :

$$E_1(\mathbb{F}_{5^2}) = \{(x, y) \in \mathbb{F} \times \mathbb{F} \mid y^2 = x^3 + x + 1\}$$

Since \mathbb{F}_{5^2} contains 25 points, in order to compute the set $E_1(\mathbb{F}_{5^2})$, we have to try $25 \cdot 25 = 625$ pairs, which is probably a bit too much for the average motivated reader. Instead, we invoke Sage to compute the curve for us. To do, we so choose the representation of \mathbb{F}_{5^2} from XXX. We get:

reference

check

check

reference

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add reference

```
sage: F5 = GF(5)
                                                                                  374
3251
    sage: F5t. < t> = F5[]
                                                                                  375
3252
    sage: P = F5t(t^2+2)
                                                                                  376
3253
    sage: P.is irreducible()
                                                                                  377
3254
    True
                                                                                  378
3255
    sage: F5_2.<t> = GF(5^2, name='t', modulus=P)
                                                                                  379
3256
    sage: E1F5 2 = EllipticCurve(F5 2,[1,1])
                                                                                  380
3257
    sage: E1F5 2.order()
                                                                                  381
3258
    27
                                                                                  382
3259
```

The curve $E_1(\mathbb{F}_{5^2})$ consist of 27 points, in contrast to curve $E_1(\mathbb{F}_5)$, which consists of 9 points. Printing the points gives the following:

```
E_{1}(\mathbb{F}_{5^{2}}) = \{ \mathscr{O}, (0,4), (0,1), (3,4), (3,1), (4,3), (4,2), (2,4), (2,1), \\ (4t+3,3t+4), (4t+3,2t+1), (3t+2,t), (3t+2,4t), \\ (2t+2,t), (2t+2,4t), (2t+1,4t+4), (2t+1,t+1), \\ (2t+3,3), (2t+3,2), (t+3,2t+4), (t+3,3t+1), \\ (3t+1,t+4), (3t+1,4t+1), (3t+3,3), (3t+3,2), (1,4t) \}
```

As we can see, curve $E_1(\mathbb{F}_5)$ sits inside curve $E(\mathbb{F}_{5^2})$, which is implied by \mathbb{F}_5 being a subfield of \mathbb{F}_{5^2} .

Full torsion groups The fundamental theorem of finite cyclic groups XXX <u>implies that every</u> prime factor *n* of a cyclic group's order defines a subgroup of the size of the prime factor. Such a subgroup is called an *n*-torsion group. We have seen many of those subgroups in the examples XXX and XXX.

When we consider elliptic curve extensions as defined in 5.23, we could ask what happens to the n-torsion groups in the extension. One might intuitively think that their extension just

This needs to be written (in Algebra)

add reference

add reference

parallels the extension of the curve. For example, when $E(\mathbb{F}_p)$ is a curve over prime field \mathbb{F}_p , with some n-torsion group \mathbb{G} and when we extend the curve to $E(\mathbb{F}_{p^m})$, then there is a bigger n-torsion group such that \mathbb{G} is a subgroup. This might make intuitive sense, as $E(\mathbb{F}_p)$ is a sub-curve of $E(\mathbb{F}_{p^m})$.

However, the actual situation is a bit more surprising than that. To see that, let \mathbb{F}_p be a prime field and let $E(\mathbb{F}_p)$ be an elliptic curve of order r, with embedding degree k and n-torsion group $E(\mathbb{F}_p)[n]$ for the same prime factor n of r. Then it can be shown that the n-torsion group $E(\mathbb{F}_p^m)[n]$ of a curve extension is equal to $E(\mathbb{F}_p)[n]$, as long as the power m is less than the embedding degree k of $E(\mathbb{F}_p)$.

However, for the prime power p^m , for any $m \ge k$, $E(\mathbb{F}_{p^m})[n]$ is strictly larger than $E(\mathbb{F}_p)[n]$ and contains $E(\mathbb{F}_p)[n]$ as a subgroup. We call the *n*-torsion group $E(\mathbb{F}_{p^k})[n]$ of the extension of E over \mathbb{F}_{p^k} the **full** *n*-**torsion group** of that elliptic curve. It can be shown that it contains n^2 many elements and consists of n+1 subgroups, one of which is $E(\mathbb{F}_p)[n]$.

So, roughly speaking, when we consider towers of curve extensions $E(\mathbb{F}_{p^m})$ ordered by the prime power m, then the n-torsion group stays constant for every level m, that is smaller than the embedding degree, while it suddenly blossoms into a larger group on level k with n+1 subgroups, and then stays like that for any level m larger than k. In other words, once the extension field is big enough to find one more point of order n (that is not defined over the base field), then we actually find all of the points in the full torsion group.

towers of curve extensions

Example 85. Consider curve $E_1(\mathbb{F}_5)$ again. We know that it contains a 3-torsion group and that the embedding degree of 3 is 2. From this we can deduce that we can find the full 3-torsion group $E_1[3]$ in the curve extension $E_1(\mathbb{F}_{5^2})$, the latter of which we computed in example 84.

check reference

Since that curve is small, in order to find the full 3-torsion, we can loop through all elements of $E_1(\mathbb{F}_{5^2})$ and check check the defining equation $[3]P = \mathcal{O}$. Invoking Sage, we compute as follows:

```
sage: INF = E1F5_2(0) # Point at infinity
                                                                               383
3293
    sage: L_E1_3 = []
                                                                               384
3294
    sage: for p in E1F5_2:
                                                                               385
3295
                if 3*p == INF:
                                                                               386
3296
                     L E1_3.append(p)
                                                                               387
3297
    sage: E1_3 = Set(L_E1_3) # Full 3-torsion set
                                                                               388
3298
```

We get the following result:

```
E_1[3] = \{ \mathscr{O}, (1,t), (1,4t), (2,1), (2,4), (2t+1,t+1), (2t+1,4t+4), (3t+1,t+4), (3t+1,4t+1) \}
```

Example 86. Consider the tiny jubjub curve from example 66. We know from example 82 that it contains a 5-torsion group and that the embedding degree of 5 is 4. This implies that we can find the full 5-torsion group $TJJ_13[5]$ in the curve extension $TJJ_13[\mathbb{F}_{13^4}]$.

check reference

To compute the full torsion, first observe that, since \mathbb{F}_{13^4} contains 28561 elements, computing $TJJ_13(\mathbb{F}_{13^4})$ means checking $28561^2=815730721$ elements. From each of these curve points P, we then have to check the equation $[5]P=\mathcal{O}$. Doing this for 815730721 is a bit too slow even on a computer.

check reference

Fortunately, Sage has a way to loop through points of a given order efficiently. The following Sage code provides a way to compute the full torsion group:

```
      3308
      sage: # define the extension field
      389

      3309
      sage: F13= GF(13) # prime field
      390

      3310
      sage: F13t.<t> = F13[] # polynomials over t
      391
```

```
sage: P = F13t(t^4+2) # irreducible polynomial of degree 4
                                                                             392
3311
    sage: P.is_irreducible()
                                                                             393
3312
                                                                             394
3313
    sage: F13_4.<t> = GF(13^4, name='t', modulus=P) # F {13^4}
                                                                             395
3314
    sage: TJJF13_4 = EllipticCurve(F13_4,[8,8]) # tiny jubjub
                                                                             396
3315
       extension
3316
    sage: # compute the full 5-torsion
                                                                             397
3317
    sage: L_TJJF13_4_5 = []
                                                                             398
3318
    sage: INF = TJJF13_4(0)
                                                                             399
3319
    sage: for P in INF.division_points(5): # [5]P == INF
3320
                                                                             400
                L_TJJF13_4_5.append(P)
                                                                             401
3321
    sage: len(L_TJJF13_4_5)
                                                                             402
3322
    25
                                                                             403
3323
    sage: TJJF13_4_5 = Set(L_TJJF13_4_5)
                                                                             404
3324
```

As expected, we get a group that contains $5^2 = 25$ elements. As it's rather tedious to write this group down, and as we don't need it in what follows, we forgo doing this. To see that the embedding degree 4 is actually the smallest prime power to find the full 5-torsion group, let's compute the 5-torsion group over of the tiny-jubjub curve of the extension field \mathbb{F}_{13^3} . We get the following:

```
sage: # define the extension field
                                                                            405
3330
    sage: P = F13t(t^3+2) # irreducible polynomial of degree 3
                                                                            406
3331
    sage: P.is_irreducible()
                                                                            407
                                                                            408
3333
    sage: F13_3.<t> = GF(13^3, name='t', modulus=P) # F_{13^3}
                                                                            409
3334
    sage: TJJF13_3 = EllipticCurve(F13_3,[8,8]) # tiny jubjub
                                                                            410
3335
       extension
3336
    sage: # compute the 5-torsion
                                                                            411
3337
    sage: L_TJJF13_3_5 = []
                                                                            412
3338
    sage: INF = TJJF13_3(0)
3339
                                                                            413
    sage: for P in INF.division_points(5): # [5]P == INF
                                                                            414
3340
               L_TJJF13_3_5.append(P)
                                                                            415
3341
    sage: len(L_TJJF13_3_5)
                                                                            416
3342
    5
                                                                            417
3343
    sage: TJJF13_3_5 = Set(L_TJJF13_3_5) # full $5$-torsion
                                                                            418
3344
```

As we can see, the 5-torsion group of tiny-jubjub over \mathbb{F}_{13^3} is equal to the 5-torsion group of tiny-jubjub over \mathbb{F}_{13} itself.

Example 87. Let's look at the curve Secp256k1. We know from example 67 that the curve is of some prime order r. Because of this, the only n-torsion group to consider is the curve itself, so the curve group is the r-torsion.

check reference

However, in order to find the full *r*-torsion of Secp256k1, we need to compute the embedding degree *k*. And as we have seen in XXX<u>it is at least not small</u>. However, we know from Fermat's little theorem (page 21 ff.) that a finite embedding degree must exist. It can be shown that it is given by the following 256-bit number:

add reference

k = 192986815395526992372618308347813175472927379845817397100860523586360249056

This means that the embedding degree is huge, which implies that the field extension \mathbb{F}_{p^k} is huge too. To understand how big \mathbb{F}_{p^k} is, recall that an element of \mathbb{F}_{p^m} can be represented as a

is "huge" a technical term?

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string $[x_0, ..., x_m]$ of m elements, each containing a number from the prime field \mathbb{F}_p . Now, in the case of Secp256k1, such a representation has k-many entries, each of them 256 bits in size. So, without any optimizations, representing such an element would need $k \cdot 256$ bits, which is too much to be represented in the observable universe.

Torsion subgroups As we have stated above, any full *n*-torsion group contains n+1 cyclic subgroups, two of which are of particular interest in pairing-based elliptic curve cryptography. To characterize these groups, we need to consider the so-called **Frobenius endomorphism** of an elliptic curve $E(\mathbb{F})$ over some finite field \mathbb{F} of characteristic p:

$$\pi: E(\mathbb{F}) \to E(\mathbb{F}): \begin{array}{ccc} (x,y) & \mapsto & (x^p, y^p) \\ \mathscr{O} & \mapsto & \mathscr{O} \end{array}$$
 (5.24)

It can be shown that π maps curve points to curve points. The first thing to note is that, in case \mathbb{F} is a prime field, the Frobenius endomorphism acts trivially, since $(x^p, y^p) = (x, y)$ on prime fields due to Fermat's little theorem (page 21 ff.). This means that the Frobenius map is more interesting over prime field extensions.

check reference

With the Frobenius map at hand, we can characterize two important subgroups of the full n-torsion. The first subgroup is the n-torsion group that already exists in the curve over the base field. In pairing-based cryptography, this group is usually written as \mathbb{G}_1 , assuming that the prime factor n in the definition is implicitly given. Since we know that the Frobenius map acts trivially on curves over the prime field, we can define \mathbb{G}_1 as follows:

$$\mathbb{G}_1[n] := \{ (x, y) \in E[n] \mid \pi(x, y) = (x, y) \}$$
 (5.25)

In more mathematical terms, this definition means that \mathbb{G}_1 is the **Eigenspace** of the Frobenius map with respect to the **Eigenvalue** 1.

It can be shown that there is another subgroup of the full n-torsion group that can be characterized by the Frobenius map. In the context of so-called type 3 pairing-based cryptography, this subgroup is usually called \mathbb{G}_2 and it is defined as follows:

$$\mathbb{G}_2[n] := \{ (x, y) \in E[n] \mid \pi(x, y) = [p](x, y) \}$$
 (5.26)

In mathematical terms, \mathbb{G}_2 is the **Eigenspace** of the Frobenius map with respect to the **Eigenvalue** p.

Notation and Symbols 9. If the prime factor n of a curve's order is clear from the context, we sometimes simply write \mathbb{G}_1 and \mathbb{G}_2 to mean $\mathbb{G}_1[n]$ and $\mathbb{G}_2[n]$, respectively.

It should be noted, however that other definitions of \mathbb{G}_2 also exists in the literature. However, in the context of pairing-based cryptography, this is the most common one. It is particularly useful because we can define hash functions that map into \mathbb{G}_2 , which is not possible for all subgroups of the full *n*-torsion.

Example 88. Consider the curve $E_1(\mathbb{F}_5)$ from example 65 again. As we have seen, this curve has the embedding degree k=2, and a full 3-torsion group is given as follows:

$$E_1[3] = \{ \mathscr{O}, (2,1), (2,4), (1,t), (1,4t), (2t+1,t+1), (2t+1,4t+4), (3t+1,t+4), (3t+1,4t+1) \}$$
(5.27)

According to the general theory, $E_1[3]$ contains 4 subgroups, and we can characterize the subgroups \mathbb{G}_1 and \mathbb{G}_2 using the Frobenius endomorphism. Unfortunately, at the time of writing,

S: either add more explanation or move to a footnote

type 3
pairingbased
cryptography

add references?

Sage does not have a predefined Frobenius endomorphism for elliptic curves, so we have to use the Frobenius endomorphism of the underlying field as a temporary workaround. We compute as follows:

```
sage: L_G1 = []
                                                                                    419
3389
     sage: for P in E1_3:
                                                                                    420
3390
                 PiP = E1F5_2([a.frobenius() for a in P]) # pi(P)
                                                                                    421
3391
                 if P == PiP:
     . . . . :
                                                                                    422
3392
                      L_G1.append(P)
                                                                                    423
     . . . . :
3393
    sage: G1 = Set(L_G1)
                                                                                    424
3394
```

As expected, the group $\mathbb{G}_1 = \{ \mathscr{O}, (2,4), (2,1) \}$ is identical to the 3-torsion group of the (unextended) curve over the prime field $E_1(\mathbb{F}_5)$. We can use almost the same algorithm to compute the group \mathbb{G}_2 and get the following:

```
sage: L_G2 = []
                                                                                     425
3398
     sage: for P in E1 3:
                                                                                     426
3399
                 PiP = E1F5_2([a.frobenius() for a in P]) # pi(P)
                                                                                     427
3400
                 pP = 5*P \# [5]P
     . . . . :
                                                                                     428
3401
                 if pP == PiP:
                                                                                     429
     . . . . :
3402
                      L_G2.append(P)
                                                                                     430
3403
    sage: G2 = Set(L_G2)
                                                                                     431
3404
```

Thus, we have computed the second subgroup of the full 3-torsion group of curve E_1 as the set $\mathbb{G}_2 = \{ \mathscr{O}, (1,t), (1,4t) \}$.

Example 89. Consider the tiny-jubjub curve TJJ_13 from example 66. In example 86, we computed its full 5 torsion, which is a group that has 6 subgroups. We compute G1 using Sage as follows:

reference

check

```
sage: L_TJJ_G1 = []
                                                                                        32 ference
3410
     sage: for P in TJJF13_4_5:
                                                                                       433
3411
                 PiP = TJJF13_4([a.frobenius() for a in P]) # pi(P)
                                                                                       434
3412
                  if P == PiP:
                                                                                       435
     . . . . :
3413
                       L TJJ G1.append(P)
                                                                                       436
3414
     sage: TJJ_G1 = Set(L_TJJ_G1)
                                                                                       437
3415
    We get \mathbb{G}1 = \{ \mathcal{O}, (7,2), (8,8), (8,5), (7,11) \}
3416
     sage: L_TJJ_G1 = []
                                                                                       438
3417
     sage: for P in TJJF13_4_5:
                                                                                       439
3418
                 PiP = TJJF13_4([a.frobenius() for a in P]) # pi(P)
                                                                                       440
3419
                 pP = 13*P # [5]P
     . . . . :
                                                                                       441
3420
                  if pP == PiP:
     . . . . :
                                                                                       442
3421
                       L_TJJ_G1.append(P)
                                                                                       443
3422
     sage: TJJ_G1 = Set(L_TJJ_G1)
                                                                                       444
```

Example 90. Consider Bitcoin's curve Secp256k1 again. Since the group \mathbb{G}_1 is identical to the torsion group of the unextended curve, and since Secp256k1 has prime order, we know that, in this case, \mathbb{G}_1 is identical to Secp256k1. It is however, infeasible not to compute not only \mathbb{G}_2 itself, but to even compute an average element of \mathbb{G}_2 , as elements need too much storage to be

 $\mathbb{G}_2 = \{ \mathscr{O}, (9t^2 + 7, t^3 + 11t_1), (9t^2 + 7, 12t^3 + 2t_1), (4t^2 + 7, 5t^3 + 10t_1), (4t^2 + 7, 8t^3 + 3t_1) \}$

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3429 representable in this universe.

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The Weil pairing In this part, we consider a pairing function defined on the subgroups $\mathbb{G}_1[r]$ and $\mathbb{G}_2[r]$ of the full r-torsion E[r] of a short Weierstraß elliptic curve. To be more precise, let $E(\mathbb{F}_p)$ be an elliptic curve of embedding degree k such that r is a prime factor of its order. Then the **Weil pairing** is a bilinear, non-degenerate map:

$$e(\cdot,\cdot): \mathbb{G}_1[r] \times \mathbb{G}_2[r] \to \mathbb{F}_{p^k}; (P,Q) \mapsto (-1)^r \cdot \frac{f_{r,P}(Q)}{f_{r,Q}(P)}$$
 (5.28)

The extension field elements $f_{r,P}(Q), f_{r,Q}(P) \in \mathbb{F}_{p^k}$ are computed by Miller's algorithm:

check floating of algorithm

```
Algorithm 7 Miller's algorithm for short Weierstraß curves y^2 = x^3 + ax + b
```

```
Require: r > 3, P \in E[r], Q \in E[r] and
    b_0, \dots, b_t \in \{0, 1\} with r = b_0 \cdot 2^0 + b_1 \cdot 2^1 + \dots + b_t \cdot 2^t and b_t = 1
    procedure MILLER'S ALGORITHM(P, Q)
          if P = \mathcal{O} or Q = \mathcal{O} or P = Q then
                 return f_{r,P}(Q) \leftarrow (-1)^r
          end if
           (x_T, y_T) \leftarrow (x_P, y_P)
          f_1 \leftarrow 1
          f_2 \leftarrow 1
          for j \leftarrow t - 1, \dots, 0 do
                m \leftarrow \frac{3 \cdot x_T^2 + a}{2 \cdot y_T}
f_1 \leftarrow f_1^2 \cdot (y_Q - y_T - m \cdot (x_Q - x_T))
f_2 \leftarrow f_2^2 \cdot (x_Q + 2x_T - m^2)
x_{2T} \leftarrow m^2 - 2x_T
                 y_{2T} \leftarrow -y_T - m \cdot (x_{2T} - x_T)
                 (x_T, y_T) \leftarrow (x_{2T}, y_{2T})
                 if b_j = 1 then
                       m \leftarrow \frac{y_T - y_P}{x_T - x_P}
                       f_1 \leftarrow f_1 \cdot (y_Q - y_T - m \cdot (x_Q - x_T))
                       f_2 \leftarrow f_2 \cdot (x_Q + (x_P + x_T) - m^2)
                       x_{T+P} \leftarrow m^2 - x_T - x_P
                       y_{T+P} \leftarrow -y_T - m \cdot (x_{T+P} - x_T)
                        (x_T, y_T) \leftarrow (x_{T+P}, y_{T+P})
                 end if
          end for
          f_1 \leftarrow f_1 \cdot (x_Q - x_T)
          return f_{r,P}(Q) \leftarrow \frac{f_1}{f_2}
    end procedure
```

Understanding how the algorithm works in detail requires the concept of **divisors**, which is outside of the scope this book. The interested reader might look at XXX.

In real-world applications of pairing-friendly elliptic curves, the embedding degree is usually a small number like 2, 4, 6 or 12, and the number r is the largest prime factor of the curve's order.

add references

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Example 91. Consider curve $E_1(\mathbb{F}_5)$ from example 65. Since the only prime factor of the check 3440 group's order is 3, we cannot compute the Weil pairing on this group using our definition of Miller's algorithm. In fact, since \mathbb{G}_1 is of order 3, executing the if statement on line XXX will lead to a "division by zero" error in the computation of the slope m. 3443

reference add reference

Example 92. Consider the tiny-jubjub curve $TJJ_13(\mathbb{F}_{13})$ from example 66 again. We want to instantiate the general definition of the Weil pairing for this example. To do so, recall that, as we have see in example 82, its embedding degree is 4, and that we have the following type-3 pairing groups (where \mathbb{G}_1 and \mathbb{G}_2 are subgroups of the full 5-torsion found in the curve $TJJ_13(\mathbb{F}_{13^4})$):

check reference

> check reference

$$\mathbb{G}_1 = \{ \mathscr{O}, (7,2), (8,8), (8,5), (7,11) \}$$

$$\mathbb{G}_2 = \{ \mathscr{O}, (9t^2 + 7, t^3 + 11t), (9t^2 + 7, 12t^3 + 2t), (4t^2 + 7, 5t^3 + 10t), (4t^2 + 7, 8t^3 + 3t) \}$$

The type-3 Weil pairing is a map $e(\cdot,\cdot):\mathbb{G}_1\times\mathbb{G}_2\to\mathbb{F}_{13^4}$. From the first if-statement in Miller's algorithm, we can deduce that $e(\mathcal{O}, Q) = 1$ as well as $e(P, \mathcal{O}) = 1$ for all arguments $P \in \mathbb{G}_1$ and $Q \in \mathbb{G}_2$. In order to compute a non-trivial Weil pairing, we choose the arguments P = (7,2) and $Q = (9t^2 + 7, 12t^3 + 2t)$.

To compute the pairing $e((7,2), (9t^2+7, 12t^3+2t))$, we have to compute the extension field elements $f_{5,P}(Q)$ and $f_{5,Q}(P)$ by applying Miller's algorithm. Do do so, observe that we have $5 = 1 \cdot 2^0 + 0 \cdot 2^1 + 1 \cdot 2^2$, so we get t = 2 as well as $b_0 = 1$, $b_1 = 0$ and $b_2 = 1$. The loop therefore needs to be executed two times.

Computing $f_{5,P}(Q)$, we initiate $(x_T, y_T) = (7,2)$ as well as $f_1 = 1$ and $f_2 = 1$. Then we proceed as follows:

$$m = \frac{3 \cdot x_T^2 + a}{2 \cdot y_T}$$
$$= \frac{3 \cdot 2^2 + 1}{2 \cdot 4} = \frac{3}{3}$$
$$= 1$$

$$f_1 = f_1^2 \cdot (y_Q - y_T - m \cdot (x_Q - x_T))$$

= $1^2 \cdot (t - 4 - 1 \cdot (1 - 2)) = t - 4 + 1$
= $t + 2$

$$f_2 = f_2^2 \cdot (x_Q + 2x_T - m^2)$$

= 1² \cdot (1 + 2 \cdot 2 - 1²) = (1 + 4 - 1)
= 4

$$x_{2T} = m^2 - 2x_T$$

= 1² - 2 \cdot 2 = -3
= 2

$$y_{2T} = -y_T - m \cdot (x_{2T} - x_T)$$

= -4 - 1 \cdot (2 - 2) = -4
= 1

We update $(x_T, y_T) = (2, 1)$ and, since $b_0 = 1$, we have to execute the if statement on line XXX in the **for** loop. However, since we only loop a single time, we don't need to compute y_{T+P} , since we only need the updated x_T in the final step. We get:

$$m = \frac{y_T - y_P}{x_T - x_P}$$
$$= \frac{1 - 4}{2 - x_P}$$

$$f_1 = f_1 \cdot (y_Q - y_T - m \cdot (x_Q - x_T))$$

$$f_2 = f_2 \cdot (x_Q + (x_P + x_T) - m^2)$$

$$x_{T+P} = m^2 - x_T - x_P$$

add reference

should all lines of all algorithms be numbered?

5.3 Hashing to Curves

Elliptic curve cryptography frequently requires the ability to hash data onto elliptic curves. If the order of the curve is not a prime number, hashing to prime number subgroups is also of

importance. In the context of pairing-friendly curves, it is also sometimes necessary to hash specifically onto the group \mathbb{G}_1 or \mathbb{G}_2 .

As we have seen in section 4.1.2, many general methods are known for hashing into groups in general, and finite cyclic groups in particular. As elliptic groups are cyclic, those methods can be utilized in this case, too. However, in what follows we want to describe some methods specific to elliptic curves that are frequently used in real-world applications.

check reference

Try-and-increment hash functions One of the most straight-forward ways of hashing a bit-string onto an elliptic curve point in a secure way is to use a cryptographic hash function together with one of the methods we described in section 4.1.2 to hash to the modular arithmetics base field of the curve. Ideally, the hash function generates an image that is at least one bit longer than the bit representation of the base field modulus.

check reference

The image in the base field can then be interpreted as the x coordinate of the curve point, and the two possible y coordinates are derived from the curve equation, while one of the bits that exceeded the modulus determines which of the two y coordinates to choose.

Such an approach would be deterministic and easy to implement, and it would conserve the cryptographic properties of the original hash function. However, not all *x* coordinates generated in such a way will result in quadratic residues when inserted into the defining equation. It follows that not all field elements give rise to actual curve points. In fact, on a prime field, only half of the field elements are quadratic residues. Hence, assuming an even distribution of the hash values in the field, this method would fail to generate a curve point in about half of the attempts.

One way to account for this problem is the so-called **try-and-increment** method. Its basic assumption is that, when hashing different values, the result will eventually lead to a valid curve point.

Therefore, instead of simply hashing a string *s* to the field, we has the concatenation of *s* with additional bytes to the field instead. In other words, we use a try-and-increment hash as described in 5 is used. If the first try of hashing to the field does not result in a valid curve point, the counter is incremented, and the hashing is repeated again. This is done until a valid curve point is found.

check reference

This method has a number of advantages: It is relatively easy to implement in code, and it maintains the cryptographic properties of the original hash function. However, it is not guaranteed to find a valid curve point, as there is a chance that all possible values in the chosen size of the counter will fail to generate a quadratic residue. Fortunately, it is possible to make the probability for this arbitrarily small by choosing large enough counters and relying on the (approximate) uniformity of the hash-to-field function.

check if the algorithm is floated properly

If the curve is not of prime order, the result will be a general curve point that might not be in the "large" prime-order subgroup. In this case, a **cofactor clearing** step is then necessary to project the curve point onto the subgroup. This is done by scalar multiplication with the cofactor of prime order with respect to the curves order.

Example 93. Consider the tiny jubjub curve from example 66. We want to construct a try-and-increment hash function that hashes a binary string s of arbitrary length onto the large prime-order subgroup of size 5.

check reference

Since the curve, as well as our targeted subgroup, is defined over the field \mathbb{F}_{13} , and the binary representation of 13 is 13.bits()=1101, we apply SHA256 from Sage's hashlib library on the concatenation s||c for some binary counter string, and use the first 4 bits of the image to try to hash into \mathbb{F}_{13} . In case we are able to hash to a value z such that $z^3 + 8 \cdot z + 8$ is a quadratic residue in \mathbb{F}_{13} , we use the 5-th bit to decide which of the two possible roots of $z^3 + 8 \cdot z + 8$ we

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Algorithm 8 Hash-to- $E(\mathbb{F}_r)$ **Require:** $r \in \mathbb{Z}$ with r.nbits() = k and $s \in \{0, 1\}^*$ **Require:** Curve equation $y^2 = x^3 + ax + b$ over \mathbb{F}_r **procedure** TRY-AND-INCREMENT(r, k, s) $c \leftarrow 0$ repeat $s' \leftarrow s || c_bits()$ $z \leftarrow H(s')_0 \cdot 20' + H(s')_1 \cdot 21' + \dots + H(s')_k \cdot 2^k$ $x \leftarrow z^3 + a \cdot z + b$ $c \leftarrow c + 1$ **until** z < r and $x^{\frac{r-1}{2}} \mod r = 1$ if $H(s')_{k+1} == 0$ then $y \leftarrow \sqrt{x} \# (\text{root in } \mathbb{F}_r)$ else $y \leftarrow r - \sqrt{x} \# (\text{root in } \mathbb{F}_r)$ end if return (x, y)end procedure **Ensure:** $(x,y) \in E(\mathbb{F}_r)$

will choose as the y coordinate. The result is a curve point different from the point at infinity. To project it to a point of \mathbb{G}_1 , we multiply it with the cofactor 4. If the result is still not the point at infinity, it is the result of the hash.

To make this concrete, let s = '10011001111010100000111' be our binary string that we want to hash onto \mathbb{G}_1 . We use a 4-bit binary counter starting at zero, that is, we choose c = 0000. Invoking Sage, we define the try-hash function as follows:

```
sage: import hashlib
                                                                                    445
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    sage: def try_hash(s,c):
3509
                                                                                    446
                 s 1 = s+c
                                                                                    447
3510
                 hasher = hashlib.sha256(s_1.encode('utf-8'))
     . . . . :
                                                                                    448
3511
                 digest = hasher.hexdigest()
                                                                                    449
     . . . . :
3512
                 d = Integer(digest, base=16)
                                                                                    450
3513
                 sign = d.str(2)[-5:-4]
     . . . . :
                                                                                    451
3514
                 d = d.str(2)[-4:]
     . . . . :
                                                                                    452
3515
                 z = Integer(d, base=2)
                                                                                    453
3516
                 return (z, sign)
                                                                                    454
3517
    sage: try_hash('10011001111010110100000111','0000')
                                                                                    455
3518
     (15, '1')
                                                                                    456
3519
```

As we can see, our first attempt to hash into \mathbb{F}_{13} was not successful, as 15 is not a number in \mathbb{F}_{13} , so we increment the binary counter by 1 and try again:

```
3522 sage: try_hash('10011001111010100000111','0001') 457
3523 (3, '0') 458
```

With this try, we found a hash into \mathbb{F}_{13} . However, this point is not guaranteed to define a curve point. To see that, we insert z=3 into the right side of the Weierstraß equation of the tiny.jubjub curve, and compute $3^3+8*3+8=7$. However, 7 is not a quadratic residue in

 \mathbb{F}_{13} , since $7^{\frac{13-1}{2}} = 7^6 = 12 = -1$. This means that 3 is a not a suitable point, and we have to increment the counter two more times:

Since $6^3 + 8 \cdot 6 + 8 = 12$, and we have $\sqrt{12} \in \{5,8\}$, we finally found the valid x coordinate x = 6 for the curve point hash. Now, since the sign bit of this hash is 1, we choose the larger root y = 8 as the y coordinate and get the following hash which is a valid curve point point on the tiny jubjub curve:

$$H('10011001111010110100000111') = (6,8)$$

In order to project this onto the "large" prime-order subgroup, we have to do cofactor clearing, that is, we have to multiply the point with the cofactor 4. We get the following:

$$[4](6,8) = \mathcal{O}$$

This means that the hash value is still not right. We therefore have to increment the counter two more times again, until we finally find a correct hash to \mathbb{G}_1 :

Since $12^3 + 8 \cdot 12 + 8 = 12$, and we have $\sqrt{12} \in \{5, 8\}$, we found another valid x coordinate x = 12 for the curve point hash. Since the sign bit of this hash is 0, we choose the smaller root y = 5 as the y coordinate, and get the following hash, which is a valid curve point point on the tiny jubjub curve:

$$H('10011001111010110100000111') = (12,5)$$

In order to project this onto the "large" prime-order subgroup we have to do cofactor clearing, again? that is, we have to multiply the point with the cofactor 4. We get the following:

$$[4](12,5) = (8,5)$$

So, hashing the binary string '10011001111010100000111' onto \mathbb{G}_1 gives the hash value (8,5) as a result.

5.4 Constructing elliptic curves

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Cryptographically secure elliptic curves like Secp256k1 from example 67 have been known for quite some time. Given the latest advancements of cryptography, however, it is often necessary to design and instantiate elliptic curves from scratch that satisfy certain very specific properties.

For example, in the context of SNARK development, it was necessary to design a curve that can be efficiently implemented inside of a so-called circuit in order to enable primitives likeliptic curve signature schemes in a zero-knowledge proof. Such a curve is given by the Babyjubjub curve (66 and we have paralleled its definition by introducing the tiny-jubjub curve from

check reference

circuit

signature schemes

this was called "pen-

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example XX. Clarify difference between baby- pen- and tiny-jubjub. As we have seen, those curves are instances of so-called twisted Edwards curves, and as such have easy to implement addition laws that work without branching. However, we introduced the tiny-jubjub curve out of thin air, as we just gave the curve parameters without explaining how we came up with them.

Another requirement in the context of many so-called pairing-based zero-knowledge proofing systems is the existence of a suitable, pairing-friendly curve with a specified security level and a low embedding degree as defined in 5.2.0.1. Famous examples are the BLS_12 and the NMT curves.

The major goal of this section is to explain the most important method of designing elliptic curves with predefined properties from scratch, called the complex multiplication method. We will apply this method in section XXX to synthesize a particular BLS 6 curve, which is one of the most insecure curves, but it will serve as the main curve to build our pen-and-paper SNARKs on. As we will see, this curve has a "large" prime factor subgroup of order 13, which implies that we can use our tiny-jubjub curve to implement certain elliptic curve cryptographic primitives in circuits over that BLS 6 curve.

Before we introduce the complex multiplication method, we have to explain a few properties of elliptic curves that are of key importance in understanding the complex multiplication method.

The Trace of Frobenius To understand the complex multiplication method of elliptic curves, we have to define the so-called **trace** of an elliptic curve first.

We know from XXX that elliptic curves over finite fields are cyclic groups of finite order. Therefore, an interesting question is whether it is possible to estimate the number of elements that this curve contains. Since an affine short Weierstraß curve consists of pairs (x, y) of elements from a finite field \mathbb{F}_q plus the point at infinity, and the field \mathbb{F}_q contains q elements, the number of curve points cannot be arbitrarily large, since it can contain at most $q^2 + 1$ many elements.

There is however, a more precise estimation, usually called the **Hasse bound**. To understand it, let $E(\mathbb{F}_q)$ be an affine short Weierstraß curve over a finite field \mathbb{F}_w of order q, and let $|E(\mathbb{F}_q)|$ be the order of the curve. Then there is an integer $t \in \mathbb{Z}$, called the **trace of Frobenius** of the curve, such that $|t| \le 2\sqrt{q}$ and the following equation holds:

$$|E(\mathbb{F})| = q + 1 - t \tag{5.29}$$

A positive trace, therefore, implies that the curve contains less points than the underlying field, whereas a negative trace means that the curve contains more points. However, the estimation $|t| \le 2\sqrt{q}$ implies that the difference is not very large in either direction, and the number of elements in an elliptic curve is always approximately in the same order of magnitude as the size of the curve's base field.

Example 94. Consider the elliptic curve $E_1(\mathbb{F}_5)$ from example 65. We know that it contains 9 curve points. Since the order of \mathbb{F}_5 is 5, we compute the trace of $E_1(\mathbb{F})$ to be t=-3, since the Hasse bound is given by the following equation:

check reference

$$9 = 5 + 1 - (-3)$$

Indeed, we have $|t| \le 2\sqrt{q}$, since $\sqrt{5} > 2.23$ and $|-3| = 3 \le 4.46 = 2 \cdot 2.23 < 2 \cdot \sqrt{5}$.

Example 95. To compute the trace of the tiny-jubjub curve, recall from example 74 that the check order of PJJ_13 is 20. Since the order of \mathbb{F}_{13} is 13, we can therefore use the Hasse bound and compute the trace as t = -6:

> 20 = 13 + 1 - (-6)(5.30)

check reference

add references

add reference

reference

text to be

written in

Algebra

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Again, we have $|t| < 2\sqrt{q}$, since $\sqrt{13} > 3.60$ and $|-6| = 6 < 7.20 = 2 \cdot 3.60 < 2 \cdot \sqrt{13}$.

Example 96. To compute the trace of Secp256k1, recall from example 67 that this curve is defined over a prime field with p elements, and that the order of that group is given by r:

check reference

p = 115792089237316195423570985008687907853269984665640564039457584007908834671663r = 115792089237316195423570985008687907852837564279074904382605163141518161494337

Using the Hesse bound r = p + 1 - t, we therefore compute t = p + 1 - r, which gives the trace of curve Secp256k1 as follows:

t = 432420386565659656852420866390673177327

As we can see, Secp256k1 contains less elements than its underlying field. However, the difference is tiny, since the order of Secp256k1 is in the same order of magnitude as the order of the underlying field. Compared to p and r, t is tiny.

```
sage: p = 1157920892373161954235709850086879078532699846656405
3592
       64039457584007908834671663
3593
    sage: r = 1157920892373161954235709850086879078528375642790749
                                                                             468
3594
       04382605163141518161494337
3595
    sage: t = p + 1 - r
                                                                              469
3596
    sage: t.nbits()
                                                                              470
3597
    129
                                                                              471
3598
    sage: abs(RR(t)) \le 2*sqrt(RR(p))
                                                                              472
3599
    True
                                                                              473
3600
```

The *j*-invariant As we have seen in XXX, two elliptic curves $E_1(\mathbb{F})$ defined by $y^2 = x^3 + ax + ax + add$ reference. b and $E_2(\mathbb{F})$ defined by $y^2 + a'x + b'$ are strictly isomorphic if and only if there is a quadratic residue $d \in \mathbb{F}$ such that $a' = ad^2$ and $b' = bd^3$.

There is, however, a more general way to classify elliptic curves over finite fields \mathbb{F}_q , based on the so-called *j*-invariant of an elliptic curve with $j(E(\mathbb{F}_q)) \in \mathbb{F}_q$, as defined below:

$$j(E(\mathbb{F}_q)) = (1728 \bmod q) \frac{4 \cdot a^3}{4 \cdot a^3 + (27 \bmod q) \cdot b^2}$$
 (5.31)

A detailed description of the *j*-invariant is beyond the scope of this book. For our present purposes, it is sufficient to note that two elliptic curves $E_1(\mathbb{F})$ and $E_2(\mathbb{F}')$ are isomorphic over the algebraic closures of \mathbb{F} and \mathbb{F}' , if and only if $\overline{\mathbb{F}} = \overline{\mathbb{F}'}$ and $j(E_1) = j(E_2)$.

algebraic

So, the *j*-invariant is an important tool to classify elliptic curves and it is needed in the complex multiplication method to decide on an actual curve instantiation that implements abstractly chosen properties.

Example 97. Consider the elliptic curve $E_1(\mathbb{F}_5)$ from example 65. We compute its *j*-invariant as follows:

check reference

$$j(E_1(\mathbb{F}_5)) = (1728 \mod 5) \frac{4 \cdot 1^3}{4 \cdot 1^3 + (27 \mod 5) \cdot 1^2}$$
$$= 3 \frac{4}{4+2}$$
$$= 3 \cdot 4$$
$$= 2$$

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Example 98. Consider the elliptic curve PJJ_13 from example 66. We compute its j-invariant check as follows:

reference

$$j(E_1(\mathbb{F}_5)) = (1728 \mod 13) \frac{4 \cdot 8^3}{4 \cdot 8^3 + (27 \mod 13) \cdot 8^2}$$

$$= 12 \cdot \frac{4 \cdot 5}{4 \cdot 5 + 1 \cdot 12}$$

$$= 12 \cdot \frac{7}{7 + 12}$$

$$= 12 \cdot 7 \cdot 6^{-1}$$

$$= 12 \cdot 7 \cdot 11$$

$$01$$

Example 99. Consider Secp256k1 from example Secp256k1. We compute its j-invariant using 3612 Sage: 3613

reference

```
sage: p = 1157920892373161954235709850086879078532699846656405
                                                                               474
3614
        64039457584007908834671663
3615
    sage: F = GF(p)
                                                                               475
3616
    sage: j = F(1728) * ((F(4) *F(0)^3) / (F(4) *F(0)^3 + F(27) *F(7)^2))
                                                                               476
3617
           j == F(0)
                                                                               477
    True
                                                                               478
3619
```

The Complex Multiplication Method As we have seen in the previous sections, elliptic curves have various defining properties, like their order, their prime factors, the embedding degree, or the cardinality (number of elements) of the base field. The complex multiplication (CM) method provides a practical way of constructing elliptic curves with pre-defined restrictions on the order and the base field.

The method usually starts by choosing a base field \mathbb{F}_q of the curve $E(\mathbb{F}_q)$ we want to construct such that $q = p^m$ for some prime number p, and " $m \in \mathbb{N}$ with $m \ge 1$. We assume p > 3to simplify things in what follows.

Next, the trace of Frobenius $t \in \mathbb{Z}$ of the curve is chosen such that p and t are coprime, that is, gcd(p,t) = 0 holds true. The choice of t also defines the curve's order r, since r = p + 1 - tby the Hasse bound (equation 5.29), so choosing t will define the large order subgroup as well as all small cofactors. r has to be defined in such a way that the elliptic curve meets the security requirements of the application it is designed for.

check reference

Note that the choice of p and t also determines the embedding degree k of any prime-order subgroup of the curve, since k is defined as the smallest number such that the prime order ndivides the number $q^k - 1$.

$$D < 0$$
 $D \mod 4 = 0 \text{ or } D \mod 4 = 1$
 $4q = t^2 + |D|v^2$
(5.32)

In order for the complex multiplication method to work, neither q nor t can be arbitrary, but must be chosen in such a way that two additional integers $D \in \mathbb{Z}$ and $v \in \mathbb{Z}$ exist and the following conditions hold:

If such numbers exist, we call D the **CM-discriminant**, and we know that we can construct a curve $E(\mathbb{F}_q)$ over a finite field \mathbb{F}_q such that the order of the curve is $|E(\mathbb{F}_q)| = q + 1 - t$.

It is the content of the complex multiplication method to actually construct such a curve, that is finding the parameters a and b from \mathbb{F}_q in the defining Weiertraß equation such that the curve has the desired order r.

Finding solutions to equation 5.29,= can be achieved in different ways, but we will forego the fine detail here. In general, it can be said that there are well-known constraints for elliptic curve families (e.g. the BLS (ECT) families) that provides families of solutions. In what follows, we will look at one type curve in the BLS-family, which gives an entire range of solutions. Are we looking at a subtype of BLS or is BLS the specific type we're referring to?

check reference

disambiguate

Assuming that the proper parameters q, t, D and v are found, we have to compute the socalled **Hilbert class polynomial** $H_D \in \mathbb{Z}[x]$ of the CM-discriminant D, which is a polynomial with integer coefficients. To do so, we first have to compute the following set:

$$ICG(D) = \{(A, B, C) \mid A, B, C \in \mathbb{Z}, D = B^2 - 4AC, gcd(A, B, C) = 1,$$

$$|B| \le A \le \sqrt{\frac{|D|}{3}}, A \le C, \text{ if } B < 0 \text{ then } |B| < A < C\}$$

One way to compute this set is to first compute the integer $A_{max} = Floor(\sqrt{\frac{|D|}{3}})$, then loop through all the integers A in the range $[0, \ldots, A_{max}]$, as well as through all the integers B in the range $[-A_{max}, \ldots, A_{max}]$, then see if there is an integer C that satisfies $D = B^2 - 4AC$ and the rest of the requirements in XXX.

add reference

To compute the Hilbert class polynomial, the so-called j-function (or j-invariant) is needed, which is a complex function defined on the upper half \mathbb{H} of the complex plane \mathbb{C} , usually written as follows:is this the same as equation eq:j-invariant1?

unify terminology

$$j: \mathbb{H} \to \mathbb{C} \tag{5.33}$$

Roughly speaking, what this means is that the *j*-functions takes complex numbers $(x+i\cdot y)$ with a positive imaginary part y>0 as inputs and returns a complex number $j(x+i\cdot y)$ as a result.

For the purposes of this book, it is not important to understand the j-function in detail, and we can use Sage to compute it in a similar way that we would use Sage to compute any other well-known function. It should be noted, however, that the computation of the j-function in Sage is sometimes prone to precision errors. For example, the j-function has a root in $\frac{-1+i\sqrt{3}}{2}$, which Sage only approximates. Therefore, when using Sage to compute the j-function, we need to take precision loss into account and possibly round to the nearest integer.

```
sage: z = ComplexField(100)(0,1)
                                                                          479
3665
    sage: z \# (0+1i)
                                                                          480
3666
    1.00000000000000000000000000000000*I
                                                                          481
3667
    sage: elliptic_j(z)
                                                                          482
    3669
                                                                          483
    sage: # j-function only defined for positive imaginary
                                                                          484
3670
       arguments
3671
    sage: z = ComplexField(100)(1,-1)
                                                                          485
3672
    sage: try:
                                                                          486
3673
               elliptic_j(z)
                                                                          487
3674
    ....: except PariError:
                                                                          488
3675
3676
               pass
                                                                          489
```

```
sage: \# root at (-1+i \text{ sqrt}(3))/2
                                                                                 490
3677
    sage: z = ComplexField(100)(-1, sqrt(3))/2
                                                                                 491
3678
    sage: elliptic_j(z)
                                                                                 492
3679
    -2.6445453750358706361219364880e-88
                                                                                 493
3680
    sage: elliptic_j(z).imag().round()
                                                                                 494
3681
                                                                                 495
3682
    sage: elliptic_j(z).real().round()
                                                                                 496
3683
                                                                                 497
```

With a way to compute the j-function and the precomputed set ICG(D) at hand, we can now compute the Hilbert class polynomial as follows:

$$H_D(x) = \Pi_{(A,B,C) \in ICG(D)} \left(x - j \left(\frac{-B + \sqrt{D}}{2A} \right) \right)$$
 (5.34)

In other words, we loop over all elements (A,B,C) from the set ICG(D) and compute the j-function at the point $\frac{-B+\sqrt{D}}{2A}$, where D is the CM-discriminant that we chose in a previous step. The result defines a factor of the Hilbert class polynomial and all factors are multiplied together.

It can be shown that the Hilbert class polynomial is an integer polynomial, but actual computations need high-precision arithmetics to avoid approximation errors that usually occur in computer approximations of the *j*-function (as shown above). So, in case the calculated Hilbert class polynomial does not have integer coefficients, we need to round the result to the nearest integer. Given that the precision we used was high enough, the result will be correct.

In the next step, we use the Hilbert class polynomial $H_D \in \mathbb{Z}[x]$, and project it to a polynomial $H_{D,q} \in \mathbb{F}_q[x]$ with coefficients in the base field \mathbb{F}_q as chosen in the first step. We do this by simply computing the new coefficients as the old coefficients modulus p, that is, if $H_D(x) = a_m x^m + a_{m-1} x^{m-1} + \ldots + a_1 x + a_0$, we compute the q-modulus of each coefficient $\tilde{a}_j = a_j \mod p$, which defines the **projected Hilbert class polynomial** as follows:

$$H_{D,p}(x) = \tilde{a}_m x^m + \tilde{a}_{m-1} x^{m-1} + \ldots + \tilde{a}_1 x + \tilde{a}_0$$

We then search for roots of $H_{D,p}$, since every root j_0 of $H_{D,p}$ defines a family of elliptic curves over \mathbb{F}_q , which all have a j-invariant 5.31 or 5.33 equal to j_0 . We can pick any root, since all of them will lead to proper curves eventually.

check reference

However, some of the curves with the correct *j*-invariant might have an order different from the one we initially decided on. Therefore, we need a way to decide on a curve with the correct order.

To compute such a curve, we have to distinguish a few different cases based on our choice of the root j_0 and of the CM-discriminant D. If $j_0 \neq 0$ or $j_0 \neq 1728 \mod q$, we compute $c_1 = \frac{j_0}{(1728 \mod q) - j_0}$, then we chose some arbitrary quadratic non-residue $c_2 \in \mathbb{F}_q$, and some arbitrary cubic non-residue $c_3 \in \mathbb{F}_q$.

The following table is guaranteed to define a curve with the correct order r = q + 1 - t for the trace of Frobenius t we initially decided on:

Definition 5.4.0.1. • Case $j_0 \neq 0$ and $j_0 \neq 1728 \mod q$. A curve with the correct order is defined by one of the following equations:

actually make this a table?

$$y^2 = x^3 + 3c_1x + 2c_1$$
 or $y^2 = x^3 + 3c_1c_2^2x + 2c_1c_2^3$ (5.35)

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• Case $j_0 = 0$ and $D \neq -3$. A curve with the correct order is defined by one of the following equations:

$$y^2 = x^3 + 1$$
 or $y^2 = x^3 + c_2^3$ (5.36)

• Case $j_0 = 0$ and D = -3. A curve with the correct order is defined by one of the following equations:

$$y^2 = x^3 + 1$$
 or $y^2 = x^3 + c_2^3$ or $y^2 = x^3 + c_3^2$ or $y^2 = x^3 + c_3^{-2}$ or $y^2 = x^3 + c_3^{-2}$ or $y^2 = x^3 + c_3^{-2}c_2^3$

• Case $j_0 = 1728 \mod q$ and $D \neq -4$. A curve with the correct order is defined by one of the following equations:

$$y^2 = x^3 + x$$
 or $y^2 = x^3 + c_2^2 x$ (5.37)

• Case $j_0 = 1728 \mod q$ and D = -4. A curve with the correct order is defined by one of the following equations:

$$y^2 = x^3 + x$$
 or $y^2 = x^3 + c_2 x$ or $y^2 = x^3 + c_2^3 x$

To decide the proper defining Weierstraß equation, we therefore have to compute the order of any of the potential curves above, and then choose the one that fits our initial requirements. Since it can be shown that the Hilbert class polynomials for the CM-discriminants D=-3 and D=-4 are given by $H_{-3,q}(x)=x$ and $H_{-4,q}=x-(1728 \bmod q)$ (EXERCISE), the previous cases are exhaustive.

exercise still to be written?

To summarize, using the complex multiplication method, it is possible to synthesize elliptic curves with predefined order over predefined base fields from scratch. However, the curves that are constructed this way are just some representatives of a larger class of curves, all of which have the same order. Therefore, in real-world applications, it is sometimes more advantageous to choose a different representative from that class. To do so recall from XXX that any curve defined by the Weierstraß equation $y^2 = x^3 + axb$ is isomorphic to a curve of the form $y^2 = x^3 + ad^2x + bd^3$ for some quadratic residue $d \in \mathbb{F}_a$.

add reference

In order to find a suitable representative (e.g. with small parameters a and b) in the last step, the curve designer might choose a quadratic residue d such that the transformed curve has the properties they wanted.

Example 100. Consider curve $E_1(\mathbb{F}_5)$ from example 65. We want to use the complex multiplication method to derive that curve from scratch. Since $E_1(\mathbb{F}_5)$ is a curve of order r = 9 over the prime field of order q = 5, we know from example 94 that its trace of Frobenius is t = -3, which also implies that q and |t| are coprime.

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We then have to find parameters $D, v \in \mathbb{Z}$ such that the criteria in 5.32 hold. We get the following:

$$4q = t^{2} + |D|v^{2} \Rightarrow$$

$$20 = (-3)^{2} + |D|v^{2} \Leftrightarrow$$

$$11 = |D|v^{2}$$

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Now, since 11 is a prime number, the only solution is |D| = 11 and v = 1 here. With D = -113733 and the Euclidean division of -11 by 4 being -11 = -3.4 + 1, we have $-11 \mod 4 = 1$, which shows that D = -11 is a proper choice. 3735

In the next step, we have to compute the Hilbert class polynomial H_{-11} . To do so, we first have to find the set ICG(D). To compute that set, observe that, since $\sqrt{\frac{|D|}{3}} \approx 1.915 < 2$, we know from $A \leq \sqrt{\frac{|D|}{3}}$ and $A \in \mathbb{Z}$ that A must be either 0 or 1.

For A = 0, we know B = 0 from the constraint $|B| \le A$. However, in this case, there could be no C satisfying $-11 = B^2 - 4AC$. So we try A = 1 and deduce $B \in \{-1, 0, 1\}$ from the constraint $|B| \le A$. The case B = -1 can be excluded, since then B < 0 has to imply |B| < A. The case B=0 can also be excluded, as there cannot be an integer C with -11=-4C, since 11 is a prime number.

This leaves the case B = 1, and we compute C = 3 from the equation $-11 = 1^2 - 4C$, which gives the solution (A, B, C) = (1, 1, 3):

$$ICG(D) = \{(1,1,3)\}$$

With the set ICG(D) at hand, we can compute the Hilbert class polynomial of D = -11. To do so, we have to insert the term $\frac{-1+\sqrt{-11}}{2\cdot 1}$ into the *j*-function. To do so, first observe that $\sqrt{-11} = i\sqrt{11}$, where i is the imaginary unit, defined by $i^2 = -1$. Using this, we can invoke Sage to compute the *j*-invariant and get the following:

$$H_{-11}(x) = x - j\left(\frac{-1 + i\sqrt{11}}{2}\right) = x + 32768$$

As we can see, in this particular case, the Hilbert class polynomial is a linear function with a single integer coefficient. In the next step, we have to project it onto a polynomial from $\mathbb{F}_5[x]$ by computing the modular 5 remainder of the coefficients 1 and 32768. We get 32768 mod 5 = 3, from which it follows that the projected Hilbert class polynomial is considered a polynomial from $\mathbb{F}_5[x]$:

$$H_{-11.5}(x) = x + 3$$

As we can see, the only root of this polynomial is j = 2, since $H_{-11,5}(2) = 2 + 3 = 0$. We therefore have a situation with $j \neq 0$ and $j \neq 1728$, which tells us that we have to compute the parameter c_1 in modular 5 arithmetics:

$$c_1 = \frac{2}{1728 - 2}$$

Since 1728 mod 5 = 3, we get $c_1 = 2$. 3744

Next, we have to check if the curve $E(\mathbb{F}_5)$ defined by the Weierstraß equation $y^2 = x^3 + 3$. $2x + 2 \cdot 2$ has the correct order. We invoke Sage, and find that the order is indeed 9, so it is a curve with the required parameters. Thus, we have successfully constructed the curve with the desired properties.

Note, however, that in real-world applications, it might be useful to choose parameters a and b that have certain properties, e.g. to be a small as possible. As we know from XXX, choosing any quadratic residue $d \in \mathbb{F}_5$ gives a curve of the same order defined by $y^2 = x^2 + ak^2x + bk^3$. Since 4 is a quadratic residue in \mathbb{F}_4 , we can transform the curve defined by $y^2 = x^3 + x + 4$ into the curve $y^2 = x^3 + 4^2 + 4 \cdot 4^3$ which gives the following:

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$$y^2 = x^3 + x + 1$$

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This is the curve $E_1(\mathbb{F}_5)$ that we used extensively throughout this book. Thus, using the complex multiplication method, we were able to derive a curve with specific properties from scratch.

Example 101. Consider the tiny jubjub curve TJJ_13 from example 66. We want to use the complex multiplication method to derive that curve from scratch. Since TJJ_13 is a curve of order r = 20 over the prime field of order q = 13, we know from example 95 that its trace of Frobenius is t = -6, which also implies that q and |t| are coprime.

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We then have to find parameters $D, v \in \mathbb{Z}$ such that 5.32 holds. We get the following:

$$4q = t^{2} + |D|v^{2} \qquad \Rightarrow$$

$$4 \cdot 13 = (-6)^{2} + |D|v^{2} \qquad \Rightarrow$$

$$52 = 36 + |D|v^{2} \qquad \Leftrightarrow$$

$$16 = |D|v^{2}$$

This equation has two solutions for (D, v), namely $(-4, \pm 2)$ and $(-16, \pm 1)$. Looking at the first solution, we know that D = -4 implies j = 1728, and the constructed curve is defined by a Weierstraß equation 5.1 that has a vanishing parameter b = 0. We can therefore conclude that choosing D = -4 will not help us reconstructing TJJ_13 . It will produce curves with order 20, just not the one we are looking for.

check reference

So we choose the second solution D=-16. In the next step, we have to compute the Hilbert class polynomial H_{-16} . To do so, we first have to find the set ICG(D). To compute that set, observe that since $\sqrt{\frac{|-16|}{3}}\approx 2.31 < 3$, we know from $A \le \sqrt{\frac{|-16|}{3}}$ and $A \in \mathbb{Z}$ that A must be in the range 0..2. So we loop through all possible values of A and through all possible values of B under the constraints $|B| \le A$, and if B < 0 then |B| < A. Then we compute potential C's from $-16 = B^2 - 4AC$. We get the following two solutions for ICG(D): we get

$$ICG(D) = \{(1,0,4), (2,0,2)\}$$

With the set ICG(D) at hand, we can compute the Hilbert class polynomial of D = -16. We can invoke Sage to compute the *j*-invariant and get the following:

$$H_{-16}(x) = \left(x - j\left(\frac{i\sqrt{16}}{2}\right)\right) \left(x - j\left(\frac{i\sqrt{16}}{4}\right)\right)$$
$$= (x - 287496)(x - 1728)$$

As we can see, in this particular case, the Hilbert class polynomial is a quadratic function with two integer coefficients. In the next step, we have to project it onto a polynomial from $\mathbb{F}_5[x]$ by computing the modular 5 remainder of the coefficients 1, 287496 and 1728. We get 287496 mod 13 = 1 and 1728 mod 13 = 2, which means that the projected Hilbert class polynomial is as follows:

$$H_{-11.5}(x) = (x-1)(x-12) = (x+12)(x+1)$$

This is considered a polynomial from $\mathbb{F}_5[x]$. Thus, we have two roots, namely j=1 and j=1 12. We already know that j=12 is the wrong root to construct the tiny jubjub curve, since 1728 mod 13 = 2, and that case is not compatible with a curve with $b \neq 0$. So we choose j=1.

Another way to decide the proper root is to compute the *j*-invariant of the tiny-jubjub curve. We get the following:

$$j(TJJ_13) = 12 \frac{4 \cdot 8^3}{4 \cdot 8^3 + 1 \cdot 8^2}$$
$$= 12 \frac{4 \cdot 5}{4 \cdot 5 + 12}$$
$$= 12 \frac{7}{7 + 12}$$
$$= 12 \frac{7}{7 + 12}$$
$$= 1$$

This is equal to the root j = 1 of the Hilbert class polynomial $H_{-16,13}$ as expected. We therefore have a situation with $j \neq 0$ and $j \neq 1728$, which tells us that we have to compute the parameter c_1 in modular 5 arithmetics:

$$c_1 = \frac{1}{12 - 1} = 6$$

Since 1728 mod 13 = 12, we get $c_1 = 6$. Then we have to check if the curve $E(\mathbb{F}_5)$ defined by the Weierstraß equation $y^2 = x^3 + 3 \cdot 6x + 2 \cdot 6$, which is equivalent to $y^2 = x^3 + 5x + 12$, has the correct order. We invoke Sage and find that the order is 8, which implies that the trace of this curve is 6, not -6 as required. So we have to consider the second possibility, and choose some quadratic non-residue $c_2 \in \mathbb{F}_{13}$. We choose $c_2 = 5$ and compute the Weierstraß equation $y^2 = x^3 + 5c_2^2 + 12c_2^3$ as follows:

$$y^2 = x^3 + 8x + 5$$

We invoke Sage and find that the order is 20, which is indeed the correct one. As we know from XXX, choosing any quadratic residue $d \in \mathbb{F}_5$ gives a curve of the same order defined by add refer $y^2 = x^2 + ad^2x + bd^3$. Since 12 is a quadratic residue in \mathbb{F}_{13} , we can transform the curve defined by $y^2 = x^3 + 8x + 5$ into the curve $y^2 = x^3 + 12^2 \cdot 8 + 5 \cdot 12^3$ which gives the following:

ence

$$y^2 = x^3 + 8x + 8$$

This is the tiny jubjub curve that we used extensively throughout this book. So using the complex multiplication method, we were able to derive a curve with specific properties from scratch.

Example 102. To consider a real-world example, we want to use the complex multiplication method in combination with Sage to compute Secp256k1 from scratch. So based on example 67, we decided to compute an elliptic curve over a prime field \mathbb{F}_p of order r for the following security parameters:

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p = 115792089237316195423570985008687907853269984665640564039457584007908834671663r = 115792089237316195423570985008687907852837564279074904382605163141518161494337

According to example 96, this gives the following trace of Frobenius:

check reference

t = 432420386565659656852420866390673177327

We also decided that we want a curve of the form $y^2 = x^3 + b$, that is, we want the parameter 3767 a to be zero. This implies that the *j*-invariant of our curve must be zero.

In a first step, we have to find a CM-discriminant D and some integer v such that the equation $4p = t^2 + |D|v^2$ is satisfied. Since we aim for a vanishing j-invariant, the first thing to try is D = -3. In this case, we can compute $v^2 = (4p - t^2)$, and if v^2 happens to be an integer that has a square root v, we are done. Invoking Sage we compute as follows:

```
sage: D = -3
                                                                               498
3773
    sage: p = 1157920892373161954235709850086879078532699846656405
                                                                               499
3774
        64039457584007908834671663
3775
    sage: r = 1157920892373161954235709850086879078528375642790749
                                                                               500
3776
        04382605163141518161494337
3777
    sage: t = p+1-r
                                                                               501
3778
    sage: v_sqr = (4*p - t^2)/abs(D)
3779
                                                                               502
    sage: v_sqr.is_integer()
                                                                               503
3780
                                                                               504
3781
    sage: v = sqrt(v_sqr)
                                                                               505
3782
    sage: v.is_integer()
                                                                               506
3783
                                                                               507
3784
    sage: 4*p == t^2 + abs(D)*v^2
                                                                               508
3785
    True
3786
                                                                               509
    sage: v
                                                                               510
3787
    303414439467246543595250775667605759171
                                                                               511
3788
```

The pair (D, v) = (-3,303414439467246543595250775667605759171) does indeed solve the equation, which tells us that there is a curve of order r over a prime field of order p, defined by a Weierstraß equation $y^2 = x^3 + b$ for some $b \in \mathbb{F}_p$. Now we need to compute b.

For D = -3, we already know that the associated Hilbert class polynomial is given by $H_{-3}(x) = x$, which gives the projected Hilbert class polynomial as $H_{-3,p} = x$ and the *j*-invariant of our curve is guaranteed to be j = 0. Now, looking at 5.4.0.1, we see that there are 6 possible cases to construct a curve with the correct order r. In order to construct the curves in question, we have to choose some arbitrary quadratic and cubic non-residue. So we loop through \mathbb{F}_p to find them, invoking Sage:

check reference

```
sage: F = GF(p)
                                                                                             512
3798
     sage: for c2 in F:
                                                                                             513
3799
                   try: # quadratic residue
                                                                                             514
3800
                        _{-} = c2.nth_root(2)
     . . . . :
                                                                                             515
3801
                   except ValueError: # quadratic non residue
                                                                                             516
3802
                        break
                                                                                             517
     . . . . :
3803
     sage: c2
                                                                                             518
3804
     3
                                                                                             519
3805
     sage: for c3 in F:
                                                                                             520
3806
                   try:
                                                                                             521
     . . . . :
3808
     . . . . :
                        _{-} = c3.nth_root(3)
                                                                                             522
                   except ValueError:
                                                                                             523
3809
     . . . . :
                        break
                                                                                             524
     . . . . :
3810
     sage: c3
                                                                                             525
3811
     2
                                                                                             526
3812
```

We found the quadratic non-residue $c_2 = 3$ and the cubic non-residue $c_3 = 2$. Using those numbers, we check the six cases against the expected order r of the curve we want to

3815 synthesize:

```
sage: C1 = EllipticCurve(F,[0,1])
                                                                                   527
3816
    sage: C1.order() == r
                                                                                   528
3817
    False
                                                                                   529
3818
    sage: C2 = EllipticCurve(F, [0, c2^3])
                                                                                   530
3819
    sage: C2.order() == r
                                                                                   531
3820
                                                                                   532
3821
    sage: C3 = EllipticCurve(F, [0, c3^2])
                                                                                   533
3822
    sage: C3.order() == r
                                                                                   534
3823
                                                                                   535
3824
    sage: C4 = EllipticCurve(F, [0, c3^2*c2^3])
3825
                                                                                   536
    sage: C4.order() == r
                                                                                   537
3826
                                                                                   538
3827
    sage: C5 = EllipticCurve(F, [0, c3^(-2)])
                                                                                   539
3828
    sage: C5.order() == r
                                                                                   540
3829
    False
                                                                                   541
3830
    sage: C6 = EllipticCurve(F, [0, c3^{(-2)}*c2^{3}])
                                                                                   542
3831
    sage: C6.order() == r
3832
                                                                                   543
                                                                                   544
3833
```

As expected, we found an elliptic curve of the correct order r over a prime field of size p. In principle, we are done, as we have found a curve with the same basic properties as Secp256k1. However, the curve is defined by the following equation, which uses a very large parameter b_1 , and so it might perform too slowly in certain algorithms.

```
y^2 = x^3 + 86844066927987146567678238756515930889952488499230423029593188005931626003754
```

It is also not very elegant to be written down by hand. It might therefore be advantageous to find an isomorphic curve with the smallest possible parameter b_2 . In order to find such a b_2 , we have to choose a quadratic residue d such that $b_2 = b_1 \cdot d^3$ is as small as possible. To do so, we rewrite the last equation into the following form:

what does this mean?

$$d = \sqrt[3]{\frac{b_2}{b_1}}$$

Then we invoke Sage to loop through values $b_2 \in \mathbb{F}_p$ until it finds some number such that the quotient $\frac{b_2}{b_1}$ has a cube root d and this cube root itself is a quadratic residue.

```
sage: b1=86844066927987146567678238756515930889952488499230423 545
3836
         029593188005931626003754
3837
     sage: for b2 in F:
                                                                                            546
3838
     . . . . :
                   try:
                                                                                            547
3839
     . . . . :
                        d = (b2/b1) .nth_root(3)
                                                                                            548
3841
     . . . . :
                        try:
                                                                                            549
                                = d.nth root(2)
                                                                                            550
3842
     . . . . :
                              if d != 0:
                                                                                            551
     . . . . :
3843
                                   break
     . . . . :
                                                                                            552
3844
                        except ValueError:
                                                                                            553
     . . . . :
3845
                              pass
                                                                                            554
     . . . . :
3846
                   except ValueError:
     . . . . :
                                                                                            555
3847
3848
     . . . . :
                        pass
                                                                                            556
```

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sage: b2 557 7 558

Indeed, the smallest possible value is $b_2 = 7$ and the defining Weierstraß equation of a curve over \mathbb{F}_p with prime order r is $y^2 = x^3 + 7$, which we might call Secp256k1. As we have just seen, the complex multiplication method is powerful enough to derive cryptographically secure curves like Secp256k1 from scratch.

The BLS6_6 pen-and-paper curve In this paragraph, we summarize our understanding of elliptic curves to derive our main pen-and-paper example for the rest of the book. To do so, we want to use the complex multiplication method to derive a pairing-friendly elliptic curve that has similar properties to curves that are used in actual cryptographic protocols. However, we design the curve specifically to be useful in pen-and-paper examples, which mostly means that the curve should contain only a few points so that we are able to derive exhaustive addition and pairing tables.

A well-understood family of pairing-friendly curves is the the group of BLS curves (STUFF ABOUT THE HISTORY AND THE NAMING CONVENTION), which are derived in [XXX]. BLS curves are particularly useful in our case if the embedding degree k satisfies $k \equiv 6 \pmod{0}$, this part Of course, the smallest embedding degree k that satisfies this congruency is k = 6 and we therefore aim for a BLS6 curve as our main pen-and-paper example.

To apply the complex multiplication method from page 110 ff., recall that this method starts with a definition of the base field \mathbb{F}_{p^m} , as well as the trace of Frobenius t and the order of the curve. If the order $p^m + 1 - t$ is not a prime number, then the order r of the largest prime factor group needs to be controlled.

In the case of BLS_6 curves, the parameter m is chosen to be 1, which means that the curves are defined over prime fields. All relevant parameters p, t and r are then themselves parameterized by the following three polynomials:

$$r(x) = \Phi_6(x)$$

$$t(x) = x + 1$$

$$q(x) = \frac{1}{3}(x - 1)^2(x^2 - x + 1) + x$$
(5.38)

In the equations above, Φ_6 is the 6-th cyclotomic polynomial and $x \in \mathbb{N}$ is a parameter that the designer has to choose in such a way that the evaluation of p, t and r at the point x^{\perp} gives integers that have the proper size to meet the security requirements of the curve that they want to design. It is then guaranteed that the complex multiplication method can be used in combination with those parameters to define an elliptic curve with CM-discriminant D=-3, embedding degree k = 6, and curve equation $y^2 = x^3 + b$ for some $b \in \mathbb{F}_p$.

For example, if the curve should target the 128-bit security level, due to the Pholaard-rho attack (TODO) the parameter r should be prime number of at least 256 bits.

In order to design the smallest BLS_6 curve, we therefore have to find a parameter x such that r(x), t(x) and q(x) are the smallest natural numbers that satisfy q(x) > 3 and r(x) > 3.

We therefore initiate the design process of our BLS6 curve by looking up the 6-th cyclotomic polynomial, which is $\Phi_6 = x^2 - x + 1$, and then insert small values for x into the defining write up

add reference

check reference

cyclotomic polynomial

Pholaardrho attack

todo

¹The smallest BLS curve will also be the most insecure BLS curve. However, since our goal with this curve is ease of pen-and-paper computation rather than security, it fits the purposes of this book.

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polynomials r, t, q. We get the following results:

$$x = 1 (r(x), t(x), q(x)) (1,2,1)$$

$$x = 2 (r(x), t(x), q(x)) (3,3,3)$$

$$x = 3 (r(x), t(x), q(x)) (7,4, \frac{37}{3})$$

$$x = 4 (r(x), t(x), q(x)) (13,5,43)$$

Since q(1) = 1 is not a prime number, the first x that gives a proper curve is x = 2. However, such a curve would be defined over a base field of characteristic 3, and we would rather like to avoid that. We therefore find x = 4, which defines a curve over the prime field of characteristic 43 that has a trace of Frobenius t = 5 and a larger order prime group of size t = 13.

vhy?

Since the prime field \mathbb{F}_{43} has 43 elements and 43's binary representation is $43_2 = 101011$, which consists of 6 digits, the name of our pen-and-paper curve should be *BLS6_6*, since its is common to name a BLS curve by its embedding degree and the bit-length of the modulus in the base field. We call *BLS6_6* the **moon-math-curve**.

Based on 5.29, we know that the Hasse bound implies that *BLS6*_6 will contain exactly 39 elements. Since the prime factorization of 39 is $39 = 3 \cdot 13$, we have a "large" prime factor group of size 13, as expected, and a small cofactor group of size 3. Fortunately, a subgroup of order 13 is well suited for our purposes, as 13 elements can be easily handled in the associated addition, scalar multiplication and pairing tables in a pen-and-paper style.

check reference

We can check that the embedding degree is indeed 6 as expected, since k = 6 is the smallest number k such that r = 13 divides $43^k - 1$.

```
      3899
      sage: for k in range(1,42): # Fermat's little theorem
      559

      3900
      if (43^k-1)%13 == 0:
      560

      3901
      break
      561

      3902
      sage: k
      562

      3903
      6
      563
```

In order to compute the defining equation $y^2 = x^3 + ax + b$ of BLS6-6, we use the complex multiplication method as described in 5.4. The goal is to find $a, b \in \mathbb{F}_{43}$ representations that are particularly nice to work with. The authors of XXX showed that the CM-discriminant of every BLS curve is D = -3 and, indeed, the following equation has the four solutions $(D, v) \in \{(-3, -7), (-3, 7), (-49, -1), (-49, 1)\}$ if D is required to be negative, as expected:

check reference

what does this mean?

add reference

$$4p = t^{2} + |D|v^{2}$$

$$4 \cdot 43 = 5^{2} + |D|v^{2}$$

$$172 = 25 + |D|v^{2}$$

$$49 = |D|v^{2}$$

$$\Rightarrow$$

This means that D = -3 is indeed a proper CM-discriminant, and we can deduce that the parameter a has to be 0, and that the Hilbert class polynomial is given by $H_{-3,43}(x) = x$.

This implies that the *j*-invariant of *BLS6*_6 is given by $j(BLS6_6) = 0$. We therefore have to look at case XXX in table 5.4.0.1 to derive a parameter *b*. To decide the proper case for $j_0 = 0$ and D = -3, we therefore have to choose some arbitrary quadratic non-residue c_2 and cubic non-residue c_3 in \mathbb{F}_{43} . We choose $c_2 = 5$ and $c_3 = 36$. We check these with Sage:

add reference

check **6e**ference

sage: F43 = GF(43)

```
sage: c2 = F43(5)
                                                                                      565
3916
     ....: try: # quadratic residue
                                                                                      566
3917
                 c2.nth_root(2)
                                                                                      567
3918
     ....: except ValueError: # quadratic non residue
                                                                                      568
3919
                 c2
     . . . . :
                                                                                      569
3920
    sage: c3 = F43(36)
                                                                                      570
3921
                                                                                      571
     ....: try:
3922
     . . . . :
                 c3.nth_root(3)
                                                                                      572
3923
     ....: except ValueError:
                                                                                      573
3924
                 с3
                                                                                      574
3925
```

Using those numbers we check the six possible cases from 5.4.0.1 against the the expected order 39 of the curve we want to synthesize:

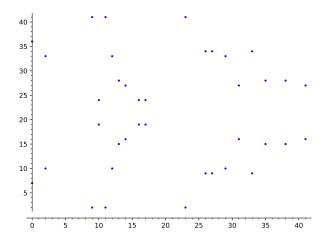
check reference

```
sage: BLS61 = EllipticCurve(F43,[0,1])
                                                                                575
3928
    sage: BLS61.order() == 39
                                                                                576
3929
    False
                                                                                577
3930
    sage: BLS62 = EllipticCurve(F43, [0, c2^3])
                                                                                578
3931
    sage: BLS62.order() == 39
                                                                                579
3932
    False
                                                                                580
3933
    sage: BLS63 = EllipticCurve(F43, [0, c3^2])
3934
                                                                                581
    sage: BLS63.order() == 39
                                                                                582
3935
    True
                                                                                583
3936
    sage: BLS64 = EllipticCurve(F43,[0,c3^2*c2^3])
                                                                                584
3937
    sage: BLS64.order() == 39
                                                                                585
3938
    False
                                                                                586
3939
    sage: BLS65 = EllipticCurve(F43,[0,c3^(-2)])
                                                                                587
3940
    sage: BLS65.order() == 39
                                                                                588
3941
    False
                                                                                589
3942
    sage: BLS66 = EllipticCurve(F43, [0, c3^{(-2)}*c2^{(3)}])
                                                                                590
3943
    sage: BLS66.order() == 39
                                                                                591
3944
    False
                                                                                592
3945
    sage: BLS6 = BLS63 # our BLS6 curve in the book
                                                                                593
```

As expected, we found an elliptic curve of the correct order 39 over a prime field of size 43, defined by the following equation:

$$BLS6_6 := \{(x, y) \mid y^2 = x^3 + 6 \text{ for all } x, y \in \mathbb{F}_{43}\}$$
 (5.39)

There are other choices for b, such as b = 10 or b = 23, but all these curves are isomorphic, and hence represent the same curve in a different way. Since BLS6-6 only contains 39 points ,it is possible to give a visual impression of the curve:



As we can see, our curve has some desirable properties: it does not contain self-inverse points, that is, points with y = 0. It follows that the addition law can be optimized, since the branch for those cases can be eliminated.

Summarizing the previous procedure, we have used the method of Barreto, Lynn and Scott to construct a pairing-friendly elliptic curve of embedding degree 6. However, in order to do elliptic curve cryptography on this curve, note that, since the order of $BLS6_6$ is 39, its group of rational points is not a finite cyclic group of prime order. We therefore have to find a suitable subgroup as our main target. Since $39 = 13 \cdot 3$, we know that the curve must contain a "large" prime-order group of size 13 and a small cofactor group of order 3.

add reference

The following step is to construct this group. One way to do so is to find a generator. We can achieve this by choosing an arbitrary element of the group that is not the point at infinity, and then multiply that point with the cofactor of the group's order. If the result is not the point at infinity, the result will be a generator. If it is the point at infinity we have to choose a different element.

In order to find a generator for the large order subgroup of size 13, we first notice that the cofactor of 13 is 3, since $39 = 3 \cdot 13$. We then need to construct an arbitrary element from $BLS6_6$. To do so in a pen-and-paper style, we can choose some $arbitraryx \in \mathbb{F}_{43}$ and see if there is some solution $y \in \mathbb{F}_{43}$ that satisfies the defining Weierstraß equation $y^2 = x^3 + 6$. We choose x = 9, and check that y = 2 is a proper solution:

$$y^{2} = x^{3} + 6 \qquad \Rightarrow \qquad 2^{2} = 9^{3} + 6 \qquad \Leftrightarrow \qquad 4 = 4$$

This implies that P = (9,2) is therefore a point on $BLS6_6$. To see if we can project this point onto a generator of the large order prime group $BLS6_6[13]$, we have to multiply P with the cofactor, that is, we have to compute [3](9,2). After some computation (EXERCISE) we get [3](9,2) = (13,15). Since this is not the point at infinity, we know that (13,15) must be a generator of $BLS6_6[13]$. The generator $g_{BLS6_6[13]}$, which we will use in pairing computations in the remainder of this book, is given as follows:

add exercise

$$g_{BLS6_6[13]} = (13, 15) (5.40)$$

Since $g_{BLS6_6[13]}$ is a generator, we can use it to construct the subgoup $BLS6_6[13]$ by repeatedly adding the generator to itself. Using Sage, we get the following:

sage:
$$P = BLS6(9,2)$$

```
sage: Q = 3*P
                                                                                      595
3976
     sage: Q.xy()
                                                                                      596
3977
     (13, 15)
                                                                                      597
3978
     sage: BLS6_13 = []
                                                                                      598
3979
     sage: for x in range(0,13): # cyclic of order 13
                                                                                      599
3980
                 P = x*0
                                                                                      600
3981
     . . . . :
                 BLS6 13.append(P)
                                                                                      601
3982
```

Repeatedly adding a generator to itself, as we just did, will generate small groups in logarithmic 3983 order with respect to the generator as, explained on page 43 ff. We therefore get the following check 3984 description of the large prime-order subgroup of *BLS6_6*:

reference

$$BLS6_6[13] = \{(13,15) \to (33,34) \to (38,15) \to (35,28) \to (26,34) \to (27,34) \to (27,9) \to (26,9) \to (35,15) \to (38,28) \to (33,9) \to (13,28) \to \mathscr{O}\} \quad (5.41)$$

Having a logarithmic description of this group is tremendously helpful in pen-and-paper computations. To see that, observe that we know from XXX that there is an exponential map from the scalar field \mathbb{F}_{13} to *BLS6_6*[13] with respect to our generator, which generates the group in logarithmic order:

add reference

$$[\cdot]_{(13,15)}: \mathbb{F}_{13} \to BLS6_6[13]; x \mapsto [x](13,15)$$

So, for example, we have $[1]_{(13,15)} = (13,15)$, $[7]_{(13,15)} = (27,9)$ and $[0]_{(13,15)} = \mathscr{O}$ and so on. The relevant point here is that we can use this representation to do computations in BLS6_6[13] efficiently in our head using XXX, as in the following example:

add reference

$$(27,34) \oplus (33,9) = [6](13,15) \oplus [11](13,15)$$
$$= [6+11](13,15)$$
$$= [4](13,15)$$
$$= (35,28)$$

So XXX is really all we need to do computations in BLS6_6[13] in this book efficiently. However, out of convenience, the following picture lists the entire addition table of that group, as it might be useful in pen-and-paper computations:

add reference

•	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)
0	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)
(13, 15)	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0
(33, 34)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27, 9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)
(38, 15)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)	(33, 34)
(35, 28)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)	(33, 34)	(38, 15)
(26, 34)	(26, 34)	(27, 34)	(27, 9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)
(27, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	Ø	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)
(27,9)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)
(26, 9)	(26,9)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)	(33,34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)
(35, 15)	(35, 15)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)
(38, 28)	(38, 28)	(33,9)	(13, 28)	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)
(33,9)	(33,9)	(13, 28)	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26,9)	(35, 15)	(38, 28)
(13, 28)	(13, 28)	0	(13, 15)	(33, 34)	(38, 15)	(35, 28)	(26, 34)	(27, 34)	(27,9)	(26, 9)	(35, 15)	(38, 28)	(33,9)

Now that we have constructed a "large" cyclic prime-order subgroup of $BLS6_6$ suitable for many pen-and-paper computations in elliptic curve cryptography, we have to look at how to do pairings in this context. We know that $BLS6_6$ is a pairing-friendly curve by design, since it has a small embedding degree k=6. It is therefore possible to compute Weil pairings efficiently. However, in order to do so, we have to decide the groups \mathbb{G}_1 and \mathbb{G}_2 as explained in exercise 73.

check reference

Since $BLS6_6$ has two non-trivial subgroups, it would be possible to use any of them as the n-torsion group. However, in cryptography, the only secure choice is to use the large prime-order subgroup, which in our case is $BLS6_6[13]$. We therefore decide to consider the 13-torsion and define $G_1[13]$ as the first argument for the Weil pairing function:

$$\mathbb{G}_{1}[13] = \{ (13,15) \to (33,34) \to (38,15) \to (35,28) \to (26,34) \to (27,34) \to (27,9) \to (26,9) \to (35,15) \to (38,28) \to (33,9) \to (13,28) \to \mathscr{O} \}$$

In order to construct the domain for the second argument, we need to construct $\mathbb{G}_2[13]$, which, according to the general theory, should be defined by those elements P of the full 13-torsion group $BLS6_6[13]$ that are mapped to $43 \cdot P$ under the Frobenius endomorphism (equation 5.24).

check reference

To compute $\mathbb{G}_2[13]$, we therefore have to find the full 13-torsion group first. To do so, we use the technique from XXX, which tells us that the full 13-torsion can be found in the curve extension over the extension field \mathbb{F}_{436} , since the embedding degree of *BLS6_6* is 6:

 $BLS6_6 := \{(x,y) \mid y^2 = x^3 + 6 \text{ for all } x, y \in \mathbb{F}_{436} \}$ (5.42)

Thus, we have to construct \mathbb{F}_{436} , a field that contains 6321363049 elements. In order to do so, we use the procedure of XXX and start by choosing a non-reducible polynomial of degree 6 from the ring of polynomials $\mathbb{F}_{43}[t]$. We choose $p(t) = t^6 + 6$. Using Sage, we get the following:

add reference

```
sage: F43 = GF(43)
                                                                                      602
4007
     sage: F43t.<t> = F43[]
                                                                                      603
4008
     sage: p = F43t(t^6+6)
                                                                                      604
4009
     sage: p.is_irreducible()
                                                                                      605
4010
     True
                                                                                      606
4011
     sage: F43 6.\langle v \rangle = GF(43^6, name='v', modulus=p)
                                                                                      607
4012
```

Recall from XXX that elements $x \in \mathbb{F}_{436}$ can be seen as polynomials $a_0 + a_1v + a_2v^2 + ... + a_5v^5$ with the usual addition of polynomials and multiplication modulo $t^6 + 6$.

add reference

In order to compute $\mathbb{G}_2[13]$, we first have to extend $BLS6_6$ to \mathbb{F}_{436} , that is, we keep the defining equation, but expand the domain from \mathbb{F}_{43} to \mathbb{F}_{436} . After that, we have to find at least one element P from that curve that is not the point at infinity, is in the full 13-torsion and satisfies the identity $\pi(P) = [43]P$. We can then use this element as our generator of $\mathbb{G}_2[13]$ and construct all other elements by repeatedly adding the generator to itself.

Since $BLS6(\mathbb{F}_{436})$ contains 6321251664 elements, it's not a good strategy to simply loop through all elements. Fortunately, Sage has a way to loop through elements from the torsion group directly:

```
sage: BLS6 = EllipticCurve (F43_6,[0 ,6]) # curve extension 608
sage: INF = BLS6(0) # point at infinity 609
```

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```
sage: for P in INF.division_points(13): # full 13-torsion
                                                                                    610
4025
           \# PI(P) == [q]P
     . . . . :
                                                                                    611
4026
                 if P.order() == 13: # exclude point at infinity
                                                                                    612
4027
                      PiP = BLS6([a.frobenius() for a in P])
                                                                                    613
     . . . . :
4028
                      qP = 43*P
                                                                                    614
     . . . . :
4029
                      if PiP == qP:
                                                                                    615
     . . . . :
4030
                           break
                                                                                    616
4031
     sage: P.xy()
                                                                                    617
4032
     (7*v^2, 16*v^3)
                                                                                    618
4033
```

We found an element from the full 13-torsion that is in the Eigenspace of the Eigenvalue 43, which implies that it is an element of $\mathbb{G}_2[13]$. As $\mathbb{G}_2[13]$ is cyclic of prime order, this element must be a generator:

$$g_{\mathbb{G}_2[13]} = (7v^2, 16v^3) \tag{5.43}$$

We can use this generator to compute \mathbb{G}_2 in logarithmic order with respect to $g_{\mathbb{G}_{[13]}}$. Using Sage we get the following:

$$\mathbb{G}_{2} = \{ (7v^{2}, 16v^{3}) \to (10v^{2}, 28v^{3}) \to (42v^{2}, 16v^{3}) \to (37v^{2}, 27v^{3}) \to (16v^{2}, 28v^{3}) \to (17v^{2}, 28v^{3}) \to (17v^{2}, 15v^{3}) \to (16v^{2}, 15v^{3}) \to (37v^{2}, 16v^{3}) \to (42v^{2}, 27v^{3}) \to (10v^{2}, 15v^{3}) \to (7v^{2}, 27v^{3}) \to \mathscr{O} \}$$

Again, having a logarithmic description of $\mathbb{G}_2[13]$ is tremendously helpful in pen-and-paper computations, as it reduces complicated computation in the extended curves to modular 13 arithmetics, as in the following example:

$$(17v^{2}, 28v^{3}) \oplus (10v^{2}, 15v^{2}) = [6](7v^{2}, 16v^{3}) \oplus [11](7v^{2}, 16v^{3})$$
$$= [6+11](7v^{2}, 16v^{3})$$
$$= [4](7v^{2}, 16v^{3})$$
$$= (37v^{2}, 27v^{3})$$

So XXX is really all we need to do computations in $\mathbb{G}_2[13]$ in this book efficiently.

To summarize the previous steps, we have found two subgroups, $\mathbb{G}_1[13]$ and $\mathbb{G}_2[13]$ suitable to do Weil pairings on *BLS6*_6 as explained in 5.28. Using the logarithmic order XXX of $\mathbb{G}_1[13]$, the logarithmic order XXX_of $\mathbb{G}_2[13]$ and the bilinearity in 5.44, we can do Weil pairings on *BLS6*_6 in a pen-and-paper style:

$$e([k_1]g_{BLS6_6[13]}, [k_2]g_{\mathbb{G}_2[13]}) = e(g_{BLS6_6[13]}, g_{\mathbb{G}_2[13]})^{k_1 \cdot k_2}$$
(5.44)

Observe that the Weil pairing between our two generators is given by the following identity:

$$e(g_{BLS6}, g_{\mathbb{G}_2[13]}) = 5v^5 + 16v^4 + 16v^3 + 15v^2 + 3v + 41$$
 (5.45)

add reference

check reference

add reference

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Hashing to pairing groups We give various constructions to hash into \mathbb{G}_1 and \mathbb{G}_2 .

We start with hashing to the scalar field... TO APPEAR

None of these techniques work for hashing into \mathbb{G}_2 . We therefore implement Pederson's Hash for BLS6.

finish writing this up

We start with \mathbb{G}_1 . Our goal is to define an 12-bit bounded hash function:

$$H_1: \{0,1\}^{12} \to \mathbb{G}_1$$

Since $12 = 3 \cdot 4$ we "randomly" select 4 uniformly distributed generators $\{(38,15),(35,28),(27,34),(38,28)\}$ from \mathbb{G}_1 and use the pseudo-random function from XXX. Therefore, we have to choose a set of 4 randomly generated invertible elements from \mathbb{F}_{13} for every generator. We choose the following:

add reference

(38,15) : {2,7,5,9} (35,28) : {11,4,7,7} (27,34) : {5,3,7,12} (38,28) : {6,5,1,8}

4059 Our hash function is then computed as follows:

$$H_1(x_{11}, x_1, \dots, x_0) = [2 \cdot 7^{x_{11}} \cdot 5^{x_{10}} \cdot 9^{x_9}](38, 15) + [11 \cdot 4^{x_8} \cdot 7^{x_7} \cdot 7^{x_6}](35, 28) + [5 \cdot 3^{x_5} \cdot 7^{x_4} \cdot 12^{x_3}](27, 34) + [6 \cdot 5^{x_2} \cdot 1^{x_1} \cdot 8^{x_0}](38, 28)$$

Note that $a^x = 1$ when x = 0. Hence, those terms can be omitted in the computation. In particular, the hash of the 12-bit zero string is given as follows:

correct computations

WRONG - ORDERING - REDO

$$H_1(0) = [2](38,15) + [11](35,28) + [5](27,34) + [6](38,28) =$$

 $(27,34) + (26,34) + (35,28) + (26,9) = (33,9) + (13,28) = (38,28)$

The hash of 011010101100 is given as follows:

fill in missing parts

$$H_1(011010101100) = WRONG - ORDERING - REDO$$

$$[2 \cdot 7^{0} \cdot 5^{1} \cdot 9^{1}](38,15) + [11 \cdot 4^{0} \cdot 7^{1} \cdot 7^{0}](35,28) + [5 \cdot 3^{1} \cdot 7^{0} \cdot 12^{1}](27,34) + [6 \cdot 5^{1} \cdot 1^{0} \cdot 8^{0}](38,28) =$$

$$[2 \cdot 5 \cdot 9](38,15) + [11 \cdot 7](35,28) + [5 \cdot 3 \cdot 12](27,34) + [6 \cdot 5](38,28) =$$

$$[12](38,15) + [12](35,28) + [11](27,34) + [4](38,28) =$$

TOAPPEAR

We can use the same technique to define a 12-bit bounded hash function in \mathbb{G}_2 :

$$H_2: \{0,1\}^{12} \rightarrow \mathbb{G}_2$$

Again, we "randomly" select 4 uniformly distributed generators $\{(7v^2, 16v^3), (42v^2, 16v^3), (17v^2, 15v^3), (10v^2, 15v^3)\}$ from \mathbb{G}_2 , and use the pseudo-random function from XXX. Therefore, we have to choose a set of 4 randomly generated invertible elements from \mathbb{F}_{13} for every generator:

add reference

$$(7v^2, 16v^3)$$
 : $\{8,4,5,7\}$
 $(42v^2, 16v^3)$: $\{12,1,3,8\}$
 $(17v^2, 15v^3)$: $\{2,3,9,11\}$
 $(10v^2, 15v^3)$: $\{3,6,9,10\}$

Our hash function is then computed like this:

$$H_1(x_{11}, x_{10}, \dots, x_0) = [8 \cdot 4^{x_{11}} \cdot 5^{x_{10}} \cdot 7^{x_9}](7v^2, 16v^3) + [12 \cdot 1^{x_8} \cdot 3^{x_7} \cdot 8^{x_6}](42v^2, 16v^3) + [2 \cdot 3^{x_5} \cdot 9^{x_4} \cdot 11^{x_3}](17v^2, 15v^3) + [3 \cdot 6^{x_2} \cdot 9^{x_1} \cdot 10^{x_0}](10v^2, 15v^3)$$

We extend this to a hash function that maps unbounded bitstrings to \mathbb{G}_2 by precomposing with an actual hash function like MD5, and feed the first 12 bits of its outcome into our previously defined hash function, with $TinyMD5_{\mathbb{G}_2}(s) = H_2(MD5(s)_0, ...MD5(s)_{11})$:

$$TinyMD5_{\mathbb{G}_2}: \{0,1\}^* \to \mathbb{G}_2$$

For example, since MD5("") =

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0xd41d8cd98f00b204e9800998ecf8427e, and the binary representation of the hexadecimal number 0x27e is 0010011111110, we compute $TinyMD5_{\mathbb{G}_7}$ of the empty string as follows:

$$TinyMD5_{\mathbb{G}_2}("") = H_2(MD5(s)_{11}, \dots MD5(s)_0) = H_2(001001111110) =$$

check equation

Chapter 6

Statements

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As we have seen in the informal introduction XXX, a SNARK is a short non-interactive argument of knowledge, where the knowledge-proof attests to the correctness of statements like "The proofer knows the prime factorization of a given number" or "The proofer knows the preimage to a given SHA2 digest value" and similar things. However, human-readable statements like these are imprecise and not very useful from a formal perspective.

In this chapter we therefore look more closely at ways to formalize statements in mathematically rigorous ways, useful for SNARK development. We start by introducing formal languages as a way to define statements properly (section 6.1). We will then look at algebraic circuits and rank-1 constraint systems as two particularly useful ways to define statements in certain formal languages (section 6.2). After that, we will have a look at fundamental building blocks of compilers that compile high-level languages to circuits and associated rank-1 constraint systems.

Proper statement design should be of high priority in the development of SNARKs, since unintended true statements can lead to potentially severe and almost undetectable security vulnerabilities in the applications of SNARKs.

Formal Languages 6.1

Formal languages provide the theoretical background in which statements can be formulated in a logically regious way and where proofing the correctness of any given statement can be realized by computing words in that language.

"rigorous

'proving"

One might argue that the understanding of formal languages is not very important in SNARK development and associated statement design, but terms from that field of research are standard jargon in many papers on zero-knowledge proofs. We therefore believe that at least some introduction to formal languages and how they fit into the picture of SNARK development is beneficial, mostly to give developers a better intuition about where all this is located in the bigger picture of the logic landscape. In addition, formal languages give a better understanding of what a formal proof for a statement actually is.

Roughly speaking, a formal language (or just language for short) is nothing but a set of words, th. Words, in turn, are strings of letters taken from some alphabet and formed according to some defining rules of the language.

To be more precise, let Σ be any set and Σ^* the set of all finite **tuples** (ordered lists) (x_1,\ldots,x_n) of elements x_i from Σ including the empty tuple $() \in \Sigma^*$. Then, a **language** L, in its most general definition, is nothing but a subset of Σ^* . In this context, the set Σ is called the **alphabet** of the language L, elements from Σ are called letters and elements from L are called words. The rules that specify which tuples from Σ^* belong to the language and which don't,

are called the **grammar** of the language. <u>S: I suggest adding an example based on English, e.g.</u> "tea" and "eat" are words of English, but "aet" and "tae" are not

Add example

If L_1 and L_2 are two formal languages over the same alphabet, we call L_1 and L_2 equivalent if there is a 1:1 correspondence between the words in L_1 and the words in L_2 . S: I'd add "In other words, two languages are equivalent if they generate the same set of words."

Add more explanation

Decision Functions Our previous definition of formal languages is very general and many subclasses of languages like **regular languages** or **context-free languages** are known in the literature. However, in the context of SNARK development, languages are commonly defined as **decision problems** where a so-called **deciding relation** $R \subset \Sigma^*$ decides whether a given tuple $x \in \Sigma^*$ is a word in the language or not. If $x \in R$ then x is a word in the associated language L_R and if $x \notin R$ then not. The relation R therefore summarizes the grammar of language L_R .

I'd delete this, too distracting

Unfortunately, in some literature on proof systems, $x \in R$ is often written as R(x), which is misleading since in general R is not a function but a relation in Σ^* . For the sake of this book, we therefore adopt a different point of view and work with what we might call a **decision function** instead:

$$R: \Sigma^* \to \{true, false\}$$
 (6.1)

Decision functions decide if a tuple $x \in \Sigma^*$ is an element of a language or not. In case a decision function is given, the associated language itself can be written as the set of all tuples that are decided by R, i.e as the set:

$$L_R := \{ x \in \Sigma^* \mid R(x) = true \}$$

$$(6.2)$$

In the context of formal languages and decision problems, a **statement** S is the claim that language L contains a word x, i.e a statement claims that there exist some $x \in L$. A constructive **proof** for statement S is given by some string $P \in \Sigma^*$ and such a proof is **verified** by checking R(P) = true. In this case, P is called an **instance** of the statement S.

While the term **language** might suggest a deeper relation to the well known **natural languages** like English, formal languages and natural languages differ in many ways. The following examples will provide some intuition about formal languages, highlighting the concepts of statements, proofs and instances:

Example 103 (Alternating Binary strings). To consider a very basic formal language with an almost trivial grammar, consider the set $\{0,1\}$ of the two letters 0 and 1 as our alphabet Σ and imply the rule that a proper word must consist of alternating binary letters of arbitrary length.

Then, the associated language L_{alt} is the set of all finite binary tuples, where a 1 must follow a 0 and vice versa. So, for example, $(1,0,1,0,1,0,1,0,1) \in L_{alt}$ is a proper word in this languages as is $(0) \in L_{alt}$ or the empty word $() \in L_{alt}$. However, the binary tuple $(1,0,1,0,1,0,1,1,1) \in \{0,1\}^*$ is not a proper word, as it violates the grammar of L_{alt} : the last3 letters are all 1. Furthermore, the tuple (0,A,0,A,0,A,0) is not a proper word, as not all its letters are not from the proper alphabet: we defined the alphabet Σ as the set $\{0,1\}$, and A is not part of that set.

Attempting to write the grammar of this language in a more formal way, we can define the following decision function:

$$R:\{0,1\}^* \to \{true,false\}\; ;\; (x_0,x_1,\ldots,x_n) \mapsto \begin{cases} true & x_{j-1} \neq x_j \text{ for all } 1 \leq j \leq n \\ false & \text{else} \end{cases}$$

We can use this function to decide if arbitrary binary tuples are words in L_{alt} or not. Some examples are given below:

binary tuples

```
• R(1,0,1) = true,
```

•
$$R(0) = true$$
,

•
$$R() = true,$$

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• but
$$R(1,1) = false$$
.

Inside our language L_{alt} , it makes sense to claim the following statement: "There exists an alternating string." One way to prove this statement constructively is by providing an actual instance, that is, finding actual alternating string like x = (1,0,1). Constructing string (1,0,1) therefore proves the statement "There exists an alternating string.", because it is easy to verify that R(1,0,1) = true.

Example 104 (Programming Language). Programming languages are a very important class of formal languages. For these languages, the alphabet is usually (a subset) of the ASCII table, and the grammar is defined by the rules of the programming language's compiler. Words, then, are nothing but properly written computer programs that the compiler accepts. The compiler can therefore be interpreted as the decision function.

To give an unusual example strange enough to highlight the point, consider the programming language Malbolge as defined in XXX. This language was specifically designed to be almost impossible to use and writing programs in this language is a difficult task. An interesting claim is therefore the statement: "There exists a computer program in Malbolge". As it turned out, proving this statement constructively, that is, by providing an actual instance of such a program, was not an easy task, as it took two years after the introduction of Malbolge to write a program that its compiler accepts. So, for two years, no one was able to prove the statement constructively.

To look at this high-level description more formally, we write $L_{Malbolge}$ for the language that uses the ASCII table as its alphabet and its words are tuples of ASCII letters that the Malbolge compiler accepts. Proving the statement "There exists a computer program in Malbolge" is then equivalent to the task of finding some word $x \in L_{Malbolge}$. The string

```
(=<'\#9] 6ZY327Uv4-QsqpMn\&+Ij"'E\%e{Ab\ w=\_:]Kw\%o44Uqp0/Q?xNvL:'H\%c\#DD2\land WV>gY;dts76qKJImZkj
```

is an example of such a proof, as it is excepted by the Malbolge compiler and is compiled to an executable binary that displays "Hello, World." (See XXX). In this example, the Malbolge compiler therefore serves as the verification process.

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Example 105 (The Empty Language). To see that not every language has even one word, consider the alphabet $\Sigma = \mathbb{Z}_6$, where \mathbb{Z}_6 is the ring of modular 6 arithmetics as derived in example 8 in chapter 3, together with the following decision function

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$$R_{\emptyset}: \mathbb{Z}_{6}^{*} \to \{true, false\}; (x_{1}, \dots, x_{n}) \mapsto \begin{cases} true & n = 4 \text{ and } x_{1} \cdot x_{1} = 2 \\ true & else \end{cases}$$

We write L_{\emptyset} for the associated language. As we can see from the multiplication table of \mathbb{Z}_6 in example 8 in chapter 3, the ring \mathbb{Z}_6 does not contain any element x such that $x^2 = 2$, which implies $R_{\emptyset}(x_1, \ldots, x_n) = false$ for all tuples $(x_1, \ldots, x_n) \in \Sigma^*$. The language therefore does not contain any words. Proving the statement "There exists a word in L_{\emptyset} " constructively by providing an instance is therefore impossible. The verification will never check any tuple.

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Example 106 (3-Factorization). We will use the following simple example repeatedly throughout this book. The task is to develop a SNARK that proves knowledge of three factors of an element from the finite field \mathbb{F}_{13} . There is nothing particularly useful about this example from an application point of view, however, in a sense, it is the most simple example that gives rise to a non trivial SNARK in some of the most common zero-knowledge proof systems.

Formalizing the high-level description, we use $\Sigma := \mathbb{F}_{13}$ as the underlying alphabet of this problem and define the language $L_{3.fac}$ to consists of those tuples of field elements from \mathbb{F}_{13} that contain exactly 4 letters w_1, w_2, w_3, w_4 which satisfy the equation $w_1 \cdot w_2 \cdot w_3 = w_4$.

So, for example, the tuple (2,12,4,5) is a word in $L_{3.fac}$, while neither (2,12,11), nor (2,12,4,7) nor (2,12,7,168) are words in $L_{3.fac}$ as they don't satisfy the grammar or are not define over the proper alphabet.

We can describe the language $L_{3.fac}$ more formally by introducing a decision function (as described in equation 6.1):

$$R_{3.fac}: \mathbb{F}_{13}^* \to \{true, false\} \; ; \; (x_1, \dots, x_n) \mapsto \begin{cases} true & n = 4 \text{ and } x_1 \cdot x_2 \cdot x_3 = x_4 \\ false & else \end{cases}$$

Having defined the language $L_{3.fac}$, it then makes sense to claim the statement "There is a word in $L_{3.fac}$ ". The way $L_{3.fac}$ is designed, this statement is equivalent to the statement "There are four elements w_1, w_2, w_3, w_4 from the finite field \mathbb{F}_{13} such that the equation $w_1 \cdot w_2 \cdot w_3 = w_4$ holds."

Proving the correctness of this statement constructively means to actually find some concrete field elements like $x_1 = 2$, $x_2 = 12$, $x_3 = 4$ and $x_4 = 5$ that satisfy the relation $R_{3.fac}$. The tuple (2,12,4,5) is therefore a constructive proof for the statement and the computation $R_{3.fac}(2,12,4,5) = true$ is a verification of that proof. In contrast, the tuple (2,12,4,7) is not a proof of the statement, since the check $R_{3.fac}(2,12,4,7) = false$ does not verify the proof.

Example 107 (Tiny JubJub Membership). In our main example, we derive a SNARK that proves a pair (x,y) of field elements from \mathbb{F}_{13} to be a point on the tiny jubjub <u>curve</u> in its **Edwards** form (see section 5.1.3).

In the first step, we define a language such that points on the tiny jubjub curve are in 1:1 correspondence with words in that language.

Since the tiny jubjub curve is an elliptic curve over the field \mathbb{F}_{13} , we choose the alphabet $\Sigma = \mathbb{F}_{13}$. In this case, the set \mathbb{F}_{13}^* consists of all finite strings of field elements from \mathbb{F}_{13} . To define the grammar, recall from 66 that a point on the tiny jubjub curve is a pair (x,y) of field elements such that $3 \cdot x^2 + y^2 = 1 + 8 \cdot x^2 \cdot y^2$. We can use this equation to derive the following decision function:

$$R_{tiny.jj}: \mathbb{F}_{13}^* \to \{true, false\}; (x_1, \dots, x_n) \mapsto \begin{cases} true & n = 2 \text{ and } 3 \cdot x_1^2 + x_2^2 = 1 + 8 \cdot x_1^2 \cdot x_2^2 \\ false & else \end{cases}$$

The associated language $L_{tiny.jj}$ is then given as the set of all strings from \mathbb{F}_{13}^* that are mapped onto true by $R_{tiny.jj}$. We get

$$L_{tiny.jj} = \{(x_1, \dots, x_n) \in \mathbb{F}_{13}^* \mid R_{tiny.jj(x_1, \dots, x_n) = true}\}$$

We can claim the statement "There is a word in $L_{tiny.jj}$ " and because $L_{tiny.jj}$ is defined by $R_{tiny.jj}$, this statement is equivalent to the claim "The tiny jubjub curve in its Edwards form has curve a point."

A constructive proof for this statement is a pair (x,y) of field elements that satisfies the Edwards equation. Example 66 therefore implies that the tuple (11,6) is a constructive proof

Are we using w and x interchange ably or is there a difference between them?

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and the computation $R_{tiny,jj}(11,6) = true$ is a proof verification. In contrast, the tuple (1,1) is not a proof of the statement, since the check $R_{tiny,jj}(1,1) = false$ does not verify the proof.

Exercise 39. Consider exercise XXX again. Define a decision function such that the associated language $L_{Exercise_XXX}$ consist precisely of all solutions to the equation 5x + 4 = 28 + 2x over \mathbb{F}_{13} . Provide a constructive proof for the claim: "There exist a word in $L_{Exercise_XXX}$ and verify the proof.

Exercise 40. Consider the modular 6 arithmetics \mathbb{Z}_6 from example 8 in chapter 3, the alphabet $\Sigma = \mathbb{Z}_6$ and the decision function

$$R_{example_8}: \Sigma^* \to \{true, false\} \; ; \; x \mapsto \begin{cases} true & x.len() = 1 \text{ and } 3 \cdot x + 3 = 0 \\ false & else \end{cases}$$

Compute all words in the associated language $L_{example_8}$, provide a constructive proof for the statement "There exist a word in $L_{example_example_8}$ " and verify the proof.

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Instance and Witness As we have seen in the previous paragraph, statements provide membership claims in formal languages, and instances serve as constructive proofs for those claims. However, in the context of **zero-knowledge** proof systems, our naive notion of constructive proofs is refined in such a way that its possible to hide parts of the proof instance and still be able to prove the statement. In this context, it is therefore necessary to split a proof into a **public part** called the **instance** and a private part called a **witness**.

To account for this separation of a proof instance into a public and a private part, our previous definition of formal languages needs a refinement in the context of zero-knowledge proof systems. Instead of a single alphabet, the refined definition considers two alphabets Σ_I and Σ_W , and a decision function defined as follows:

$$R: \Sigma_I^* \times \Sigma_W^* \to \{true, false\} \; ; \; (i;w) \mapsto R(i;w)$$
 (6.3)

Words are therefore tuples $(i; w) \in \Sigma_I^* \times \Sigma_W^*$ with R(i; w) = true. The refined definition differentiates between public inputs $i \in \Sigma_I$ and private inputs $w \in \Sigma_W$. The public input i is called an **instance** and the private input w is called a **witness** of R.

If a decision function is given, the associated language is defined as the set of all tuples from the underlying alphabet that are verified by the decision function:

$$L_R := \{ (i; w) \in \Sigma_I^* \times \Sigma_W^* \mid R(i; w) = true \}$$
(6.4)

In this refined context, a **statement** S is a claim that, given an instance $i \in \Sigma_I^*$, there is a witness $w \in \Sigma_W^*$ such that language L contains a word (i; w). A constructive **proof** for statement S is given by some string $P = (i; w) \in \Sigma_I^* \times \Sigma_W^*$ and a proof is **verified** by checking R(P) = true.

It is worth understanding the difference between statements as defined in XXX_and the refined notion of statements from this paragraph. While statements in the sense of the previous paragraph can be seen as membership claims, statements in the refined definition can be seen as knowledge-proofs, where a prover claims knowledge of a witness for a given instance. For a more detailed discussion on this topic see [XXX_sec 1.4]

add reference

Example 108 (SHA256 – Knowlege of Preimage). One of the most common examples in the context of zero-knowledge proof systems is the **knowledge-of-a-preimage proof** for some cryptographic hash function like *SHA*256, where a publicly known *SHA*256 digest value is

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given, and the task is to prove knowledge of a preimage for that digest under the *SHA*256 function, without revealing that preimage.

To understand this problem in detail, we have to introduce a language able to describe the knowledge-of-preimage problem in such a way that the claim "Given digest i, there is a preimage w such that SHA256(w) = i" becomes a statement in that language. Since SHA256 is a function

$$SHA256: \{0,1\}^* \rightarrow \{0,1\}^{256}$$

that maps binary strings of arbitrary length onto binary strings of length 256 and we want to prove knowledge of preimages, we have to consider binary strings of size 256 as instances and binary strings of arbitrary length as witnesses.

An appropriate alphabet Σ_I for the set of all instances and an appropriate alphabet Σ_W for the set of all witnesses is therefore given by the set $\{0,1\}$ of the two binary letters and a proper decision function is given by:

$$R_{SHA256}: \{0,1\}^* \times \{0,1\}^* \rightarrow \{true, false\};$$

$$(i;w) \mapsto \begin{cases} true & i.len() = 256, i = SHA256(w) \\ false & else \end{cases}$$

We write L_{SHA256} for the associated language and note that it consists of words, which are tuples (i; w) such that the instance i is the SHA256 image of the witness w.

Given some instance $i \in \{0,1\}^{256}$, a statement in L_{SHA256} is the claim "Given digest i, there is a preimage w such that SHA256(w) = i", which is exactly what the knowledge-of-preimage problem is about. A constructive proof for this statement is therefore given by a preimage w to the digest i and proof verification is achieved by checking that SHA256(w) = i.

Example 109 (3-factorization). To give an intuition about the implication of refined languages, consider $L_{3.fac}$ from example 106 again. As we have seen, a constructive proof in $L_{3.fac}$ is given by 4 field elements x_1 , x_2 , x_3 and x_4 from \mathbb{F}_{13} such that the product in modular 13 arithmetics of the first three elements is equal to the 4'th element.

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Splitting words from $L_{3.fac}$ into private and public parts, we can reformulate the problem and introduce different levels of privacy into the problem. For example, we could reformulate the membership statement of $L_{3.fac}$ into a statement where all factors x_1, x_2, x_3 of x_4 are private and only the product x_4 is public. A statement for this reformulation is then expressed by the claim: "Given a publicly known field element x_4 , there are three private factors of x_4 ". Assuming some instance x_4 , a constructive proof for the associated knowledge claim is then provided by any tuple (x_1, x_2, x_3) such that $x_1 \cdot x_2 \cdot x_3 = x_4$.

At this point, it is important to note that, while constructive proofs in the refinement don't look very different from constructive proofs in the original language, we will see in XXX that there are proof systems able to prove the statement (at least with high probability) without revealing anything about the factors x_1 , x_2 , or x_3 . This is why the importance of the refinement only becomes clear once more elaborate proofing methods beyond naive constructive proofs are provided.

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We can formalize this new language, which we might call $L_{3.fac_zk}$, by defining the following decision function:

Definition 6.1.0.1.

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$$R_{3.fac_zk} : \mathbb{F}_{13}^* \times \mathbb{F}_{13}^* \to \{true, false\};$$

$$((i_1, \dots, i_n); (w_1, \dots, w_m)) \mapsto \begin{cases} true & n = 1, \ m = 3, \ i_1 = w_1 \cdot w_2 \cdot w_3 \\ false & else \end{cases}$$

The associated language $L_{3.fac_zk}$ is defined by all tuples from $\mathbb{F}_{13}^* \times \mathbb{F}_{13}^*$ that are mapped onto *true* under the decision function $R_{3.fac_zk}$.

Considering the distinction we made between the instance and the witness part in $L_{3,fac}$ zk, one might ask why we chose the factors x_1 , x_2 and x_3 to be the witness and the product x_4 to be the instance and why we didn't choose another combination? This was an arbitrary choice in the example. Every other combination of private and public factors would be equally valid. For example, it would be possible to declare all variables as private or to declare all variables as public. Actual choices are determined by the application only.

Example 110 (The Tiny JubJub Curve). Consider the language $L_{tiny.jj}$ from example 107. As check we have seen, a constructive proof in $L_{tiny,ij}$ is given by a pair (x_1,x_2) of field elements from \mathbb{F}_{13} such that the pair is a point of the tiny jubjub curve in its Edwards representation.

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We look at a reasonable splitting of words from $L_{tiny,ij}$ into private and public parts. The two obvious choices are to either choose both coordinates x_1 as x_2 as public inputs, or to choose both coordinates x_1 as x_2 as private inputs.

In case both coordinates are public, we define the grammar of the associated language by introducing the following decision function:

$$R_{tiny.jj.1}: \mathbb{F}_{13}^* \times \mathbb{F}_{13}^* \to \{true, false\};$$

$$(I_1, \dots, I_n; W_1, \dots, W_m) \mapsto \begin{cases} true & n = 2, m = 0 \text{ and } 3 \cdot I_1^2 + I_2^2 = 1 + 8 \cdot I_1^2 \cdot I_2^2 \\ false & else \end{cases}$$

The language $L_{tiny,jj,1}$ is defined as the set of all strings from $\mathbb{F}_{13}^* \times \mathbb{F}_{13}^*$ that are mapped onto 4275 *true* by $R_{tiny,jj,1}$. 4276

In case both coordinates are private, we define the grammar of the associated refined language by introducing the following decision function:

$$R_{tiny.jj_zk} : \mathbb{F}_{13}^* \times \mathbb{F}_{13}^* \to \{true, false\} ;$$

$$(I_1, \dots, I_n; W_1, \dots, W_m) \mapsto \begin{cases} true & n = 0, m = m \text{ and } 3 \cdot W_1^2 + W_2^2 = 1 + 8 \cdot W_1^2 \cdot W_2^2 \\ false & else \end{cases}$$

The language L_{tiny,jj_zk} is defined as the set of all strings from $\mathbb{F}_{13}^* \times \mathbb{F}_{13}^*$ that are mapped onto true by $R_{tiny, jj}$ zk. 4278

Exercise 41. Consider the modular 6 arithmetics \mathbb{Z}_6 from example 8 in chapter 3 as alphabets Σ_I and Σ_W and the following decision function

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$$R_{linear}: \Sigma^* \times \Sigma^* \to \{true, false\};$$

$$(i; w) \mapsto \begin{cases} true & i.len() = 3 \text{ and } w.len() = 1 \text{ and } i_1 \cdot w_1 + i_2 = i_3 \\ false & else \end{cases}$$

Which of the following instances (i_1, i_2, i_3) has a proof of knowledge in L_{linear} ?

- (3,3,0)
- (2,1,0)

• (4,4,2)

Exercise 42 (Edwards Addition on Tiny JubJub). Consider the tiny-jubjub curve together with its Edwards addition law from example XXX. Define an instance alphabet Σ_I , a witness alphabet Σ_W and a decision function R_{add} with associated language L_{add} such that a string $(i; w) \in \Sigma_I^* \times \Sigma_W^*$ is a word in L_{add} if and only if i is a pair of curve points on the tiny-jubjub curve in Edwards form and w is the Edwards sum of those curve points.

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Choose some instance $i \in \Sigma_I^*$, provide a constructive proof for the statement "There is a witness $w \in \Sigma_W^*$ such that (i;w) is a word in L_{add} " and verify that proof. Then find some instance $i \in \Sigma_I^*$ such that i has no knowledge proof in L_{add} .

Modularity From a developers perspective, it is often useful to construct complex statements and their representing languages from simple ones. In the context of zero-knowledge proof systems, those simple building blocks are often called **gadgets**, and gadget libraries usually contain representations of atomic types like booleans, integers, various hash functions, elliptic curve cryptography and many more. In order to synthesize statements, developers then combine predefined gadgets into complex logic. We call the ability to combine statements into more complex statements **modularity**.

To understand the concept of modularity on the level of formal languages defined by decision functions, we need to look at the **intersection** of two languages, which exists whenever both languages are defined over the same alphabet. In this case, the intersection is a language that consists of strings which are words in both languages.

To be more precise, let L_1 and L_2 be two languages defined over the same instance and witness alphabets Σ_I and Σ_W . Then the intersection $L_1 \cap L_2$ of L_1 and L_2 is defined as

$$L_1 \cap L_2 := \{ x \mid x \in L_1 \text{ and } x \in L_2 \}$$
 (6.5)

If both languages are defined by decision functions R_1 and R_2 , the following function is a decision function for the intersection language $L_1 \cap L_2$:

$$R_{L_1 \cap L_2}: \Sigma_I^* \times \Sigma_W^* \to \{true, false\}; (i, w) \mapsto R_1(i, w) \text{ and } R_2(i, w)$$
 (6.6)

Thus, the intersection of two decision-function-based languages is a also decision-function-based language. This is important from an implementations point of view: It allows us to construct complex decision functions, their languages and associated statements from simple building blocks. Given a publicly known instance $i \in \Sigma_I^*$ a statement in an intersection language then claims knowledge of a witness that satisfies all relations simultaneously.

6.2 Statement Representations

As we have seen in the previous section, formal languages and their definitions by decision functions are a powerful tool to describe statements in a formally rigorous manner.

However, from the perspective of existing zero-knowledge proof systems, not all ways to actually represent decision functions are equally useful. Depending on the proof system, some are more suitable than others. In this section, will describe two of the most common ways to represent decision functions and their statements.

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Rank-1 Quadratic Constraint Systems 6.2.1

Although decision functions are expressible in various ways, many contemporary proofing systems require the deciding relation to be expressed in terms of a system of quadratic equations over a finite field. This is true in particular for pairing-based proofing systems like XXX, roughly because it is possible to check solutions to those equations "in the exponent" of pairingfriendly cryptographic groups.

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In this section, we will therefore have a closer look at a particular type of quadratic equation called rank-1 quadratic constraints systems, which are a common standard in zeroknowledge proof systems. We will start with a general introduction to those systems and then look at their relation to formal languages. We will look into a common way to compute solutions to those systems, and then show how a simple compiler might derive rank-1 constraint systems from more high-level programming code.

R1CS representation To understand what rank-1 (quadratic) constraint systems are in detail, let \mathbb{F} be a field, n, m and $k \in \mathbb{N}$ three numbers and a_j^i , b_j^i and $c_j^i \in \mathbb{F}$ constants from \mathbb{F} for every index $0 \le j \le n + m$ and $1 \le i \le k$. Then a rank-1 constraint system (R1CS) is defined as follows:

Definition 6.2.1.1. R1CS representation

If a rank-1 constraint system is given, the parameter k is called the **number of constraints** If a tuple $(I_1, \ldots, I_n; W_1, \ldots, W_m)$ of field elements satisfies these equations, (I_1, \ldots, I_n) is called an **instance** and (W_1, \ldots, W_m) is called an associated witness of the system.

Remark 1 (Matrix notation). The presentation of rank-1 constraint systems can be simplified using the notation of vectors and matrices, which abstracts over the indices. In fact if x = $(1,I,W) \in \mathbb{F}^{1+n+m}$ is a (n+m+1)-dimensional vector, A,B,C are $(n+m+1) \times k$ -dimensional matrices and ⊙ is the Schur/Hadamard product, then a R1CS can be written as

Schur/Hadam product

$$Ax \odot Bx = Cx$$

However, since we did not introduced matrix calculus in the book, we use XXX as the defining equation for rank-1 constraints systems. We only highlighted the matrix notation, because it is sometimes used in the literature.

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Generally speaking, the idea of a rank-1 constraint system is to keep track of all the values that any variable can assume during a computation and to bind the relationships among all those variables that are implied by the computation itself. Enforcing relations between all the steps of a computer program, the execution is then constrained to be computed in exactly the expected way without any opportunity for deviations. In this sense, solutions to rank-1 constraint systems are proofs of proper progam execution.

Example 111 (3-Factorization). To provide a better intuition of rank-1 constraint systems, con-4346 sider the language $L_{3.fac_zk}$ from example 106 again. As we have seen, $L_{3.fac_zk}$ consists of 4347 words $(I_1; W_1, W_2, W_3)$ over the alphabet \mathbb{F}_{13} such that $I_1 = W_1 \cdot W_2 \cdot W_3$. We show how to rewrite reference 4348 the decision function as a rank-1 constraint system. 4349

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Since R1CS are systems of quadratic equations, expressions like $W_1 \cdot W_2 \cdot W_3$ which contain products of more than two factors (which are therefore not quadratic) have to be rewritten in a process often called **flattening**. To flatten the defining equation $I_1 = W_1 \cdot W_2 \cdot W_3$ of $L_{3,fac,zk}$ we introcuce a new variable W_4 , which captures two of the three multiplications in $W_1 \cdot W_2 \cdot W_3$. We get the following two constraints

$$W_1 \cdot W_2 = W_4$$
 constraint 1
 $W_4 \cdot W_3 = I_1$ constraint 2

Given some instance I_1 , any solution (W_1, W_2, W_3, W_4) to this system of equations provides a solution to the original equation $I_1 = W_1 \cdot W_2 \cdot W_3$ and vice versa. Both equations are therefore equivalent in the sense that solutions are in a 1:1 correspondence.

Looking at both equations, we see how each constraint enforces a step in the computation. In fact, the first constraint forces any computation to multiply the witness W_1 and W_2 first. Otherwise it would not be possible to compute the witness W_4 , which is needed to solve the second constraint. Witness W_4 therefore expresses the constraining of an intermediate computational

At this point, one might ask why equation 1 constrains the system to compute $W_1 \cdot W_2$ first, since computing $W_2 \cdot W_3$, or $W_1 \cdot W_3$ in the beginning and then multiplying with the remaining factor gives the exact same result. The reason is that the way we designed the R1CS prohibits any of these alternative computations, which shows that R1CS are in general not unique descriptions of a language: many different R1CS are able to describe the same problem.

To see that the two quadratic equations qualify as a rank-1 constraint system, choose the parameter n = 1, m = 4 and k = 2 as well as

With this choice, the rank-1 constraint system of our 3-factorization problem can be written in its most general form as follows:

$$(a_0^1 + a_1^1 I_1 + a_2^1 W_1 + a_3^1 W_2 + a_4^1 W_3 + a_5^1 W_4) \cdot (b_0^1 + b_1^1 I_1 + b_2^1 W_1 + b_3^1 W_2 + b_4^1 W_3 + b_5^1 W_4) = (c_0^1 + c_1^1 I_1 + c_2^1 W_1 + c_3^1 W_2 + c_4^1 W_3 + c_5^1 W_4)$$

$$(a_0^2 + a_1^2 I_1 + a_2^2 W_2 + a_3^2 W_2 + a_4^2 W_3 + a_5^2 W_4) \cdot (b_0^2 + b_1^2 I_1 + b_2^2 W_2 + b_3^2 W_3 + b_5^2 W_4) = (c_0^2 + c_1^2 I_1 + c_2^2 W_2 + c_4^2 W_3 + c_5^2 W_4)$$

Example 112 (The Tiny Jubjub curve). Consider the languages $L_{tiny,jj,1}$ from example 107, check which consist of words (I_1, I_2) over the alphabet \mathbb{F}_{13} such that $3 \cdot I_1^2 + I_2^2 = 1 + 8 \cdot I_1^2 \cdot I_2^2$.

reference

We derive a rank-1 constraint system such that its associated language is equivalent to $L_{tiny, jj, 1}$. To achieve this, we first rewrite the defining equation:

$$3 \cdot I_1^2 + I_2^2 = 1 + 8 \cdot I_1^2 \cdot I_2^2 \qquad \Leftrightarrow 0 = 1 + 8 \cdot I_1^2 \cdot I_2^2 - 3 \cdot I_1^2 - I_2^2 \qquad \Leftrightarrow 0 = 1 + 8 \cdot I_1^2 \cdot I_2^2 + 10 \cdot I_1^2 + 12 \cdot I_2^2$$

Since R1CSs are systems of quadratic equations, we have to reformulate this expression into a system of quadratic equations. To do so, we have to introduce new variables that constrain intermediate steps in the computation and we have to decide if those variables should be public or private. We decide to declare all new variables as private and get the following constraints

$$I_1 \cdot I_1 = W_1$$
 constraint 1
 $I_2 \cdot I_2 = W_2$ constraint 2
 $(8 \cdot W_1) \cdot W_2 = W_3$ constraint 3
 $(12 \cdot W_2 + W_3 + 10 \cdot W_1 + 1) \cdot 1 = 0$ constraint 4

To see that these four quadratic equations qualify as a rank-1 constraint system according to definition XXX, choose the parameter n = 2, m = 3 and k = 4:

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With this choice, the rank-1 constraint system of our tiny-jubjub curve point problem can be written in its most general form as follows:

```
 \begin{pmatrix} a_0^1 + a_1^1 I_1 + a_2^1 I_2 + a_3^1 W_1 + a_4^1 W_2 + a_5^1 W_3 \end{pmatrix} \cdot \begin{pmatrix} b_0^1 + b_1^1 I_1 + b_2^1 I_2 + b_3^1 W_1 + b_4^1 W_2 + b_5^1 W_3 \end{pmatrix} = \begin{pmatrix} c_0^1 + c_1^1 I_1 + c_2^1 I_2 + c_3^1 W_1 + c_4^1 W_2 + c_5^1 W_3 \end{pmatrix} \\ \begin{pmatrix} a_0^2 + a_1^2 I_1 + a_2^2 I_2 + a_3^2 W_1 + a_4^2 W_2 + a_5^2 W_3 \end{pmatrix} \cdot \begin{pmatrix} b_0^2 + b_1^2 I_1 + b_2^2 I_2 + b_3^2 W_1 + b_4^2 W_2 + b_5^2 W_3 \end{pmatrix} = \begin{pmatrix} c_0^2 + c_1^2 I_1 + c_2^2 I_2 + c_3^2 W_1 + c_4^2 W_2 + c_5^2 W_3 \end{pmatrix} \\ \begin{pmatrix} a_0^3 + a_1^3 I_1 + a_2^3 I_2 + a_3^3 W_1 + a_4^3 W_2 + a_5^3 W_3 \end{pmatrix} \cdot \begin{pmatrix} b_0^3 + b_1^3 I_1 + b_2^3 I_2 + b_3^3 W_1 + b_4^3 W_2 + b_5^3 W_3 \end{pmatrix} = \begin{pmatrix} c_0^3 + c_1^3 I_1 + c_2^3 I_2 + c_3^3 W_1 + c_4^3 W_2 + c_5^3 W_3 \end{pmatrix} \\ \begin{pmatrix} a_0^4 + a_1^4 I_1 + a_2^4 I_2 + a_3^4 W_1 + a_4^4 W_2 + a_5^4 W_3 \end{pmatrix} \cdot \begin{pmatrix} b_0^4 + b_1^4 I_1 + b_2^4 I_2 + b_3^4 W_1 + b_4^4 W_2 + b_5^4 W_3 \end{pmatrix} = \begin{pmatrix} c_0^4 + c_1^4 I_1 + c_2^4 I_2 + c_3^4 W_1 + c_4^4 W_2 + c_5^4 W_3 \end{pmatrix}
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In what follows, we write L_{jubjub} for the associated language that consists of solutions to the R1CS.

To see that L_{jubjub} is equivalent to $L_{tiny.jj.1}$, let $(I_1,I_2;W_1,W_2,W_3)$ be a word in L_{jubjub} , then (I_1,I_2) is a word in $L_{tiny.jj.1}$, since the defining R1CS of L_{jubjub} implies that I_1 and I_2 satisfy the Edwards equation of the tiny jubjub curve. On the other hand, let (I_1,I_2) be a word in $L_{tiny.jj.1}$. Then $(I_1,I_2;I_1^2,I_2^2,8\cdot I_1^2\cdot I_2^2)$ is a word in L_{jubjub} and both maps are inverses of each other.

Exercise 43. Consider the language $L_{tiny.jj_zk}$ and define a rank-1 constraint relation with a decision function such that the associated language is equivalent to $L_{tiny.jj_zk}$.

R1CS Satisfiability To understand how rank-1 constraint systems define formal languages, observe that every R1CS over a field $\mathbb F$ defines a decision function over the alphabet $\Sigma_I \times \Sigma_W =$

 $\mathbb{F} \times \mathbb{F}$ in the following way: 4375

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$$R_{R1CS}: \mathbb{F}^* \times \mathbb{F}^* \to \{true, false\}; (I; W) \mapsto \begin{cases} true & (I; W) \text{ satisfies R1CS} \\ false & else \end{cases}$$
 (6.7)

Every R1CS therefore defines a formal language. The grammar of this language is encoded in the constraints, words are solutions to the equations and a **statement** is a knowledge claim "Given instance I, there is a witness W such that (I; W) is a solution to the rank-1 constraint system". A constructive proof to this claim is therefore an assignment of a field element to every witness variable, which is verified whenever the set of all instance and witness variables solves the R1CS.

Remark 2 (R1CS satisfiability). It should be noted that in our definition, every R1CS defines its own language. However, in more theoretical approaches, another language usually called R1CS satisfiability is often considered, which is useful when it comes to more abstract problems like expressiveness or the computational complexity of the class of all R1CS. From our perspective, the R1CS satisfiability language is obtained by the union of all R1CS languages that are in our definition. To be more precise, let the alphabet $\Sigma = \mathbb{F}$ be a field. Then

$$L_{R1CS_SAT(\mathbb{F})} = \{(i; w) \in \Sigma^* \times \Sigma^* \mid \text{there is a R1CS } R \text{ such that } R(i; w) = true\}$$

Example 113 (3-Factorization). Consider the language $L_{3,fac\ zk}$ from example 106 and the R1CS defined in example ex:3-factorization-r1cs. As we have seen in ex:3-factorization-r1cs, solutions to the R1CS are in 1:1 correspondence with solutions to the decision function of $L_{3,fac\ zk}$. Both languages are therefore equivalent in the sense that there is a 1:1 correspondence between words in both languages.

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To give an intuition of what constructive proofs in $L_{3.fac_zk}$ look like, consider the instance check $I_1 = 11$. To prove the statement "There exists a witness W such that $(I_1; W)$ is a word in $L_{3.fac_zk}$ " constructively, a proof has to provide assignments to all witness variables W_1 , W_2 , W_3 and W_4 . Since the alphabet is \mathbb{F}_{13} , an example assignment is given by W=(2,3,4,6) since $(I_1; W)$ satisfies the R1CS

$$W_1 \cdot W_2 = W_4$$
 # $2 \cdot 3 = 6$
 $W_4 \cdot W_3 = I_1$ # $6 \cdot 4 = 11$

A proper constructive proof is therefore given by P = (2,3,4,6). Of course, P is not the only 4387 possible proof for this statement. Since factorization is not unique in a field in general, another 4388 constructive proof is given by P' = (3, 5, 12, 2). 4389

Example 114 (The tiny jubjub curve). Consider the language L_{jubjub} from example 107 and check its associated R1CS. To see how constructive proofs in $L_{jub\,jub}$ look like, consider the instance $(I_1,I_2)=(11,6)$. To prove the statement "There exists a witness W such that $(I_1,I_2;W)$ is a word in L_{jubjub} " constructively, a proof has to provide assignments to all witness variables W_1 , W_2 and W_3 . Since the alphabet is \mathbb{F}_{13} , an example assignment is given by W=(4,10,8) since $(I_1, I_2; W)$ satisfies the R1CS

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$$I_1 \cdot I_1 = W_1$$
 $11 \cdot 11 = 4$ $I_2 \cdot I_2 = W_2$ $6 \cdot 6 = 10$ $(8 \cdot W_1) \cdot W_2 = W_3$ $(8 \cdot 4) \cdot 10 = 8$ $(12 \cdot W_2 + W_3 + 10 \cdot W_1 + 1) \cdot 1 = 0$ $12 \cdot 10 + 8 + 10 \cdot 4 + 1 = 0$

A proper constructive proof is therefore given by P = (4, 10, 8), which shows that the instance 4390 (11,6) is a point on the tiny jubjub curve.

Modularity As we discussed on page 135 XXX, it is often useful to construct complex statements and their representing languages from simple ones. Rank-1 constraint systems are particularly useful for this, as the intersection of two R1CS over the same alphabet results in a new R1CS over that same alphabet.

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To be more precise, let S_1 and S_2 be two R1CS over \mathbb{F} , then a new R1CS S_3 is obtained by the intersection $S_3 = S_1 \cap S_2$ of S_1 and S_2 . In this context, intersection means that both the equations of S_1 and the equations of S_2 have to be satisfied in order to provide a solution for the system S_3 .

As a consequence, developers are able to construct complex R1CS from simple ones and this modularity provides the theoretical foundation for many R1CS compilers, as we will see in XXX.

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6.2.2 Algebraic Circuits

As we have seen in the previous paragraphs, rank-1 constraint systems are quadratic equations such that solutions are knowledge proofs for the existence of words in associated languages. From the perspective of a proofer, it is therefore important to solve those equations efficiently.

However, in contrast to systems of linear equation, no general methods are known that solve systems of quadratic equations efficiently. Rank-1 constraint systems are therefore impractical from a proofers perspective and auxiliary information is needed that helps to compute solutions efficiently.

Methods which compute R1CS solutions are sometimes called **witness generator functions**. To provide a common example, we introduce another class of decision functions called **algebraic circuits**. As we will see, every algebraic circuit defines an associated R1CS and also provides an efficient way to compute solutions for that R1CS.

It can be shown that every space- and time-bounded computation is expressible as an algebraic circuit. Transforming high-level computer programs into those circuits is a process often called **flattening**.

To understand this in more detail, we will introduce our model for algebraic circuits and look at the concept of circuit execution and valid assignments. After that, we will show how to derive rank-1 constraint systems from circuits and how circuits are useful to compute solutions to their R1CS efficiently.

Algebraic circuit representation To see what algebraic circuits are, let \mathbb{F} be a field. An algebraic circuit is then a directed acyclic (multi)graph that computes a polynomial function over \mathbb{F} . Nodes with only outgoing edges (source nodes) represent the variables and constants of the function and nodes with only incoming edges (sink nodes) represent the outcome of the function. All other nodes have exactly two incoming edges and represent the defining field operations addition as well as multiplication. Graph edges represent the flow of the computation along the nodes.

To be more precise, we call a directed acyclic multi-graph $C(\mathbb{F})$ an **algebraic circuit** over \mathbb{F} in this book if the following conditions hold:

Definition 6.2.2.1. Algebraic circuit

- The set of edges has a total order.
- Every source node has a label that represents either a variable or a constant from the field
 \mathbb{F}.

- Every sink node has exactly one incoming edge and a label that represents either a variable or a constant from the field \mathbb{F} .
- Every node that is neither a source nor a sink has exactly two incoming edges and a label from the set $\{+,*\}$ that represents either addition or multiplication in \mathbb{F} .
- All outgoing edges from a node have the same label.
- Outgoing edges from a node with a label that represents a variable have a label.
- Outgoing edges from a node with a label that represents multiplication have a label, if there is at least one labeled edge in both input path.
- All incoming edges to sink nodes have a label.
- If an edge has two labels S_i and S_j it gets a new label $S_i = S_j$.
- No other edge has a label.
- Incoming edges to sink nodes that are labeled with a constant $c \in \mathbb{F}$ are labeled with the same constant. Every other edge label is taken from the set $\{W, I\}$ and indexed compatible with the order of the edge set.

It should be noted that the details in the definitions of algebraic circuits vary between different sources. We use this definition as it is conceptually straightforward and well-suited for pen-and-paper computations.

To get a better intuition of our definition, let $C(\mathbb{F})$ be an algebraic circuit. Source nodes are the inputs to the circuit and either represent variables or constants. In a similar way, sink nodes represent termination points of the circuit and are either output variables or constants. Constant sink nodes enforce computational outputs to take on certain values.

Nodes that are neither source nodes nor sink nodes are called **arithmetic gates**. Arithmetic gates that are decorated with the "+"-label are called **addition-gates** and arithmetic gates that are decorated with the "·"-label are called **multiplication-gates**. Every arithmetic gate has exactly two inputs, represented by the two incoming edges.

Since the set of edges is ordered, we can write it as $\{E_1, E_2, ..., E_n\}$ for some $n \in \mathbb{N}$ and we use those indices to index the edge labels, too. Edge labels are therefore either constants or symbols like I_j , W_j or S_j , where j is an index compatible with the edge order. Labels I_j represent instance variables, labels W_j witness variables. Labels on the outgoing edges of input variables constrain the associated variable to that edge. Every other edge defines a constraining equation in the associated R1CS. We will explain this in more detail in XXX.

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Notation and Symbols 10. In synthesizing algebraic circuits, assigning instance I_j or witness W_j labels to appropriate edges is often the final step. It is therefore convenient to not distinguish these two types of edges in previous steps. To account for that, we often simply write S_j for an edge label, indicating that the private/public property of the label is unspecified and it might represent an instance or a witness label.

Example 115 (Generalized factorization SNARK). To give a simple example of an algebraic circuit, consider our 3-factorization problem from example 106 again. To express the problem in the algebraic circuit model, consider the following function

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$$f_{3.fac}: \mathbb{F}_{13} \times \mathbb{F}_{13} \times \mathbb{F}_{13} \to \mathbb{F}_{13}; (x_1, x_2, x_3) \mapsto x_1 \cdot x_2 \cdot x_3$$

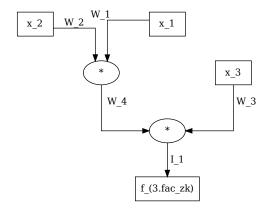
Using this function, we can describe the zero-knowledge 3-factorization problem from 106, check 4474 in the following way: Given instance $I_1 \in \mathbb{F}_{13}$, a valid witness is a preimage of $f_{3,fac}$ at the point I_1 , i.e., a valid witness consists of three values W_1 , W_2 and W_3 from \mathbb{F}_{13} such that $f_{3,fac}(W_1,W_2,W_3)=I_1.$ 4477

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To see how this function can be transformed into an algebraic circuit over \mathbb{F}_{13} , it is a common first step to introduce brackets into the function's definition and then write the operations as binary operators, in order to highlight how exactly every field operation acts on its two inputs. Due to the associativity laws in a field, we have several choices. We choose

$$f_{3.fac}(x_1, x_2, x_3) = x_1 \cdot x_2 \cdot x_3$$
 # bracket choice
= $(x_1 \cdot x_2) \cdot x_3$ # operator notation
= $MUL(MUL(x_1, x_2), x_3)$

Using this expression, we can write an associated algebraic circuit by first constraining the 4478 variables to edge labels $W_1 = x_1$, $W_2 = x_2$ and $W_3 = x_3$ as well as $I_1 = f_{3,fac}(x_1, x_2, x_3)$, taking the distinction between private and public inputs into account. We then rewrite the operator representation of $f_{3.fac}$ into circuit nodes and get the following: 4481



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In this case, the directed acyclic multi-graph is a binary tree with three leaves (the source nodes) labeled by x_1 , x_2 and x_3 , one root (the single sink node) labeled by $f(x_1, x_2, x_3)$ and two internal nodes, which are labeled as multiplication gates.

The order we use to label the edges is chosen to make the edge labeling consistent with the choice of W_4 as defined in definition 6.2.2.1. This order can be obtained by a depth-first right-to-left-first traversal algorithm.

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Example 116. To give a more realistic example of an algebraic circuit, look at the defining equation of the tiny-jubjub curve (66) again. A pair of field elements $(x,y) \in \mathbb{F}_{13}^2$ is a curve point, precisely if the following equation holds:

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$$3 \cdot x^2 + y^2 = 1 + 8 \cdot x^2 \cdot y^2$$

To understand how one might transform this identity into an algebraic circuit, we first rewrite this equation by shifting all terms to the right. We get the following:

$$3 \cdot x^{2} + y^{2} = 1 + 8 \cdot x^{2} \cdot y^{2} \qquad \Leftrightarrow 0 = 1 + 8 \cdot x^{2} \cdot y^{2} - 3 \cdot x^{2} - y^{2} \qquad \Leftrightarrow 0 = 1 + 8 \cdot x^{2} \cdot y^{2} + 10 \cdot x^{2} + 12 \cdot y^{2}$$

Then we use this expression to define a function such that all points of the tiny-jubjub curve are characterized as the function preimages at 0.

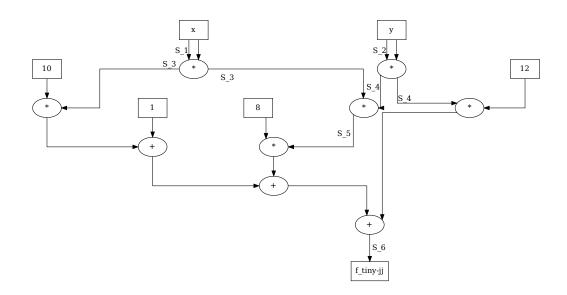
$$f_{tiny-ij}: \mathbb{F}_{13} \times \mathbb{F}_{13} \to \mathbb{F}_{13}; (x,y) \mapsto 1 + 8 \cdot x^2 \cdot y^2 + 10 \cdot x^2 + 12 \cdot y^2$$

Every pair of points $(x,y) \in \mathbb{F}^2_{13}$ with $f_{tiny-jj}(x,y) = 0$ is a point on the tiny-jubjub curve, and there are no other curve points. The preimage $f_{tiny-jj}^{-1}(0)$ is therefore a complete description of the tiny-jubjub curve.

We can transform this function into an algebraic circuit over \mathbb{F}_{13} . We first introduce brackets into potentially ambiguous expressions and then rewrite the function in terms of binary operators. We get the following:

$$f_{tiny-jj}(x,y) = 1 + 8 \cdot x^{2} \cdot y^{2} + 10 \cdot x^{2} + 12y^{2} \qquad \Leftrightarrow \\ = ((8 \cdot ((x \cdot x) \cdot (y \cdot y))) + (1 + 10 \cdot (x \cdot x))) + (12 \cdot (y \cdot y)) \qquad \Leftrightarrow \\ = ADD(ADD(MUL(8,MUL(MUL(x,x),MUL(y,y))),ADD(1,MUL(10,MUL(x,x)))),MUL(12,MUL(y,y)))$$

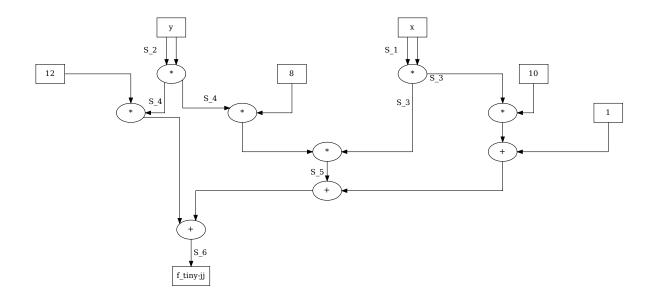
Since we haven't decided which part of the computation should be public and which part should be private, we use the unspecified symbol S to represent edge labels. Constraining all variables to edge labels $S_1 = x$, $S_2 = y$ and $S_6 = f_{tiny-jj}$, we get the following circuit, representing the function $f_{tiny-jj}$, by inductively replacing binary operators with their associated arithmetic gates:



This circuit is not a graph, but a multigraph, since there is more than one edge between some of the nodes.

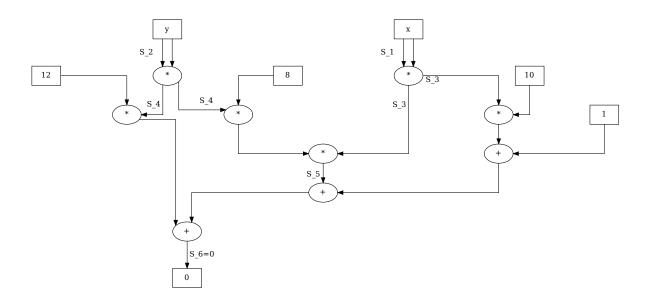
In the process of designing of circuits from functions, it should be noted that circuit representations are not unique in general. In case of the function $f_{tiny-jj}$, the circuit shape is dependent on our choice of bracketing in XXX. An alernative design i,s for example, given by the following circuit, which occurs when the bracketed expression $8 \cdot ((x \cdot x) \cdot (y \cdot y))$ is replaced by the expression $(x \cdot x) \cdot (8 \cdot (y \cdot y))$.

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Of course, both circuits represent the same function, due to the associativity and commutativity laws that hold true in any field.

With a circuit that represents the function $f_{tiny-jj}$, we can now proceed to derive a circuit that constrains arbitrary pairs (x,y) of field elements to be points on the tiny-jubjub curve. To do so, we have to constrain the output to be zero, that is, we have to constrain $S_6 = 0$. To indicate this in the circuit, we replace the output variable by the constant 0 and constrain the related edge label accordingly. We get the following:



The previous circuit enforces input values assigned to the labels S_1 and S_2 to be points on the tiny jubjub curve. However, it does not specify which labels are considered public and which are considered private. The following circuit defines the inputs to be public, while all other labels are private:



It can be shown that every space- and time-bounded computation can be transformed into an algebraic circuit. We call any process that transforms a bounded computation into a circuit **flattening**.

Circuit Execution Algebraic circuits are directed, acyclic multi-graphs, where nodes represent variables, constants, or addition and multiplication gates. In particular, every algebraic circuit with n input nodes decorated with variables symbols and m output nodes decorated with variables can be seen a function that transforms an input tuple (x_1, \ldots, x_n) from \mathbb{F}^n into an output tuple (f_1, \ldots, f_m) from \mathbb{F}^m . The transformation is done by sending values associated to nodes along their outgoing edges to other nodes. If those nodes are gates, then the values are transformed according to the gate label and the process is repeated along all edges until a sink node is reached. We call this computation **circuit execution**.

When executing a circuit, it is possible to not only compute the output values of the circuit but to derive field elements for all edges, and, in particular, for all edge labels in the circuit. The result is a tuple (S_1, S_2, \ldots, S_n) of field elements associated to all labeled edges, which we call a **valid assignment** to the circuit. In contrast, any assignment $(S'_1, S'_2, \ldots, S'_n)$ of field elements to edge labels that can not arise from circuit execution is called an **invalid assignment**.

Valid assignments can be interpreted as **proofs for proper circuit execution** because they keep a record of the computational result as well as intermediate computational steps.

Example 117 (3-factorization). Consider the 3-factorization problem from example 106 and its representation as an algebraic circuit from XXX. We know that the set of edge labels is given by $S := \{I_1; W_1, W_2, W_3, W_4\}$.

by $S := \{I_1; W_1, W_2, W_3, W_4\}$.

To understand how this circuit is executed, consider the variables $x_1 = 2$, $x_2 = 3$ as well as $x_3 = 4$. Following all edges in the graph, we get the assignments $W_1 = 2$, $W_2 = 3$ and $W_3 = 4$.

Then the assignments of W_1 and W_2 enter a multiplication gate and the output of the gate is $2 \cdot 3 = 6$, which we assign to W_4 , i.e. $W_4 = 6$. The values W_4 and W_3 then enter the second

multiplication gate and the output of the gate is $6 \cdot 4 = 11$, which we assign to I_1 , i.e. $I_1 = 11$. A valid assignment to the 3-factorization circuit $C_{3,fac}(\mathbb{F}_{13})$ is therefore given by the fol-

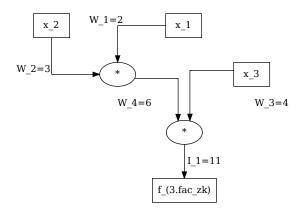
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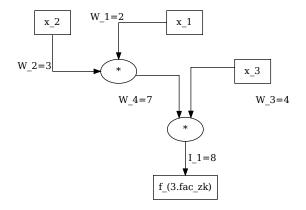
$$S_{valid} := \{11; 2, 3, 4, 6\} \tag{6.8}$$

We can visualise this assignment in the circuit as follows:



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To see what an invalid assignment looks like, consider the assignment $S_{err} := \{8; 2, 3, 4, 7\}$. In this assignment, the input values are the same as in the previous case. The associated circuit is:



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This assignment is invalid, as the assignments of I_1 and W_4 cannot be obtained by executing the circuit.

Example 118. To compute a more realistic algebraic circuit execution, consider the defining circuit $C_{tiny-jj}(\mathbb{F}_{13})$ from example 114 again. We already know from the way this circuit is constructed that any valid assignment with $S_1 = x$, $S_2 = y$ and $S_6 = 0$ will ensure that the pair (x,y) is a point on the tiny jubjub curve in its Edwards representation (equation 5.20.

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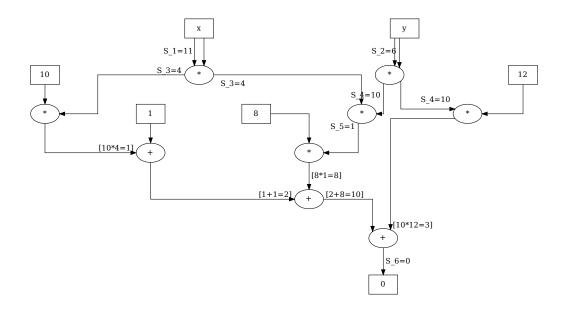
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From example 114, we know that the pair (11,6) is a proper point on the tiny-jubjub curve and we use this point as input to a circuit execution. We get the following:

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Executing the circuit, we indeed compute $S_6 = 0$ as expected, which proves that (11,6) is a point on the tiny-jubjub curve in its Edwards representation. A valid assignment of $C_{tiny-jj}(\mathbb{F}_{13})$ is therefore given by the following equation:

$$S_{tiny-jj} = \{S_1, S_2, S_3, S_4, S_5, S_6\} = \{11, 6, 4, 10, 1, 0\}$$

Circuit Satisfiability To understand how algebraic circuits give rise to formal languages, observe that every algebraic circuit $C(\mathbb{F})$ over a fields \mathbb{F} defines a decision function over the alphabet $\Sigma_I \times \Sigma_W = \mathbb{F} \times \mathbb{F}$ in the following way:

$$R_{C(\mathbb{F})}: \mathbb{F}^* \times \mathbb{F}^* \to \{true, false\}; (I; W) \mapsto \begin{cases} true & (I; W) \text{is valid assignment to } C(\mathbb{F}) \\ false & else \end{cases}$$
 (6.9)

Every algebraic circuit therefore defines a formal language. The grammar of this language is encoded in the shape of the circuit, words are assignments to edge labels that are derived from circuit execution, and **statements** are knowledge claims "Given instance I, there is a witness W such that (I;W) is a valid assignment to the circuit". A constructive proof to this claim is therefore an assignment of a field element to every witness variable, which is verified by executing the circuit to see if the assignment of the execution meets the assignment of the proof.

In the context of zero-knowledge proof systems, executing circuits is also often called **witness generation**, since in applications the instance part is usually public, while its the task of a proofer to compute the witness part.

Remark 3 (Circuit satisfiability). It should be noted that, in our definition, every circuit defines its own language. However, in more theoretical approaches another language usually called **circuit satisfiability** is often considered, which is useful when it comes to more abstract problems like expressiveness, or computational complexity of the class of **all** algebraic circuits over a given field. From our perspective the circuit satisfiability language is obtained by union of all circuit languages that are in our definition. To be more precise, let the alphabet $\Sigma = \mathbb{F}$ be a field. Then

Should we refer to R1CS satisfiability (p. 139 here?

 $L_{CIRCUIT_SAT(\mathbb{F})} = \{(i; w) \in \Sigma^* \times \Sigma^* \mid \text{there is a circuit } C(\mathbb{F}) \text{ such that } (i; w) \text{ is valid assignment} \}$

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Example 119 (3-Factorization). Consider the circuit $C_{3,fac}$ from equation 6.8 again. We call the check associated language $L_{3.fac_circ}$.

reference

To understand how a constructive proof of a statement in $L_{3.fac_circ}$ looks like, consider the instance $I_1 = 11$. To provide a proof for the statement "There exist a witness W such that $(I_1; W)$ is a word in $L_{3.fac_circ}$ " a proof therefore has to consists of proper values for the variables W_1 , W_2 , W_3 and W_4 . Any proofer therefore has to find input values for W_1 , W_2 and W_3 and then execute the circuit to compute W_4 under the assumption $I_1 = 11$.

Example XXX implies that (2,3,4,6) is a proper constructive proof and in order to verify the proof a verifier needs to execute the circuit with instance $I_1 = 11$ and inputs $W_1 = 2$, $W_2 = 3$ and $W_3 = 4$ to decide whether the proof is a valid assignment or not.

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Associated Constraint Systems As we have seen in 6.2.1.1, rank-1 constraint systems define check a way to represent statements in terms of a system of quadratic equations over finite fields, suitable for pairing-based zero-knowledge proof systems. However, those equations provide no practical way for a proofer to actually compute a solution. On the other hand, algebraic circuits can be executed in order to derive valid assignments efficiently.

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In this paragraph, we show how to transform any algebraic circuit into a rank-1 constraint system such that valid circuit assignments are in 1:1 correspondence with solutions to the associated R1CS.

To see this, let $C(\mathbb{F})$ be an algebraic circuit over a finite field \mathbb{F} , with a set of edge labels $\{S_1, S_2, \dots, S_n\}$. Then one of the following steps is executed for every edge label S_i from that

• If the edge label S_j is an outgoing edge of a multiplication gate, the Definition 6.2.2.2. R1CS gets a new quadratic constraint

$$(left input) \cdot (right input) = S_j \tag{6.10}$$

where (left input) respectively (right input) is the output from the symbolic execution of the subgraph that consists of the left respectively right input edge of this gate, and all edges and nodes that have this edge in their path, starting with constant inputs or labeled outgoing edges of other nodes.

• If the edge label S_j is an outgoing edge of an addition gate, the R1CS gets a new quadratic constraint

$$(left input + right input) \cdot 1 = S_i$$
 (6.11)

where (left input) respectively (right input) is the output from the symbolic execution of the subgraph that consists of the left respectively right input edge of this gate and all edges and nodes that have this edge in their path, starting with constant inputs or labeled outgoing edges of other nodes.

• No other edge label adds a constraint to the system.

The result of this method is a rank-1 constraint system, and, in this sense, every algebraic circuit $C(\mathbb{F})$ generates a R1CS R, which we call the **associated R1CS** of the circuit. It can be shown that a tuple of field elements (S_1, S_2, \dots, S_n) is a valid assignment to a circuit if and only if the same tuple is a solution to the associated R1CS. Circuit executions therefore compute solutions to rank-1 constraints systems efficiently.

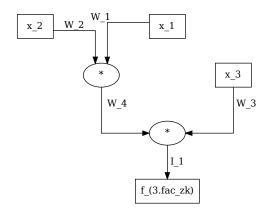
To understand the contribution of algebraic gates to the number of constraints, note that, by definition, multiplication gates have labels on their outgoing edges if and only if there is at least

one labeled edge in both input paths, or if the outgoing edge is an input to a sink node. This implies that multiplication with a constant is essentially free in the sense that it doesn't add a new constraint to the system, as long as that multiplication gate is not am input to an output

Moreover, addition gates have labels on their outgoing edges if and only if they are inputs to sink nodes. This implies that addition is essentially free in the sense that it doesn't add a new constraint to the system, as long as that addition gate is not an input to an output node.

Example 120 (3-factorization). Consider our 3-factorization problem from equation 6.8 and the check associated circuit $C_{3.fac}(\mathbb{F}_{13})$. Out task is to transform this circuit into an equivalent rank-1 constraint system.

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We start with an empty R1CS, and, in order to generate all constraints, we have to iterate over 4632 the set of edge labels $\{I_1; W_1, W_2, W_3, W_4\}$. 4633

Starting with the edge label I_1 , we see that it is an outgoing edge of a multiplication gate, and, since both input edges are labeled, we have to add the following constraint to the system:

$$(\text{left input}) \cdot (\text{right input}) = I_1 \Leftrightarrow W_4 \cdot W_3 = I_1$$

Next, we consider the edge label W_1 and, since, it's not an outgoing edge of a multiplication or 4634 addition label, we don't add a constraint to the system. The same holds true for the labels W_2 4635 and W_3 . 4636

For edge label W_4 , we see that it is an outgoing edge of a multiplication gate, and, since both input edges are labeled, we have to add the following constraint to the system:

$$(\text{left input}) \cdot (\text{right input}) = W_4 \Leftrightarrow W_2 \cdot W_1 = W_4$$

Since there are no more labeled edges, all constraints are generated, and we have to combine them to get the associated R1CS of $C_{3.fac}(\mathbb{F}_{13})$:

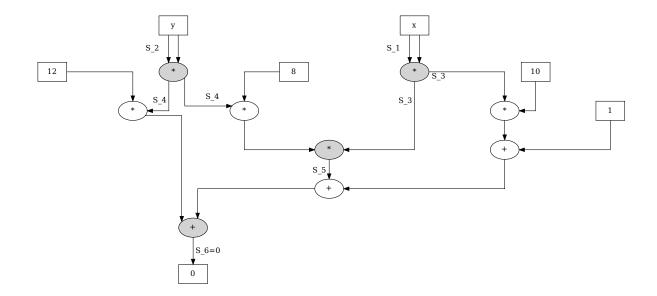
$$W_4 \cdot W_3 = I_1$$

$$W_2 \cdot W_1 = W_4$$

This system is equivalent to the R1CS we derived in example 111. The languages $L_{3,fac,zk}$ and check 4637 $L_{3.fac\ circ}$ are therefore equivalent and both the circuit as well as the R1CS are just two different reference 4638 ways of expressing the same language.

Example 121. To consider a more general transformation, we consider the tiny-jubjub circuit from example 114 again. A proper circuit is given by

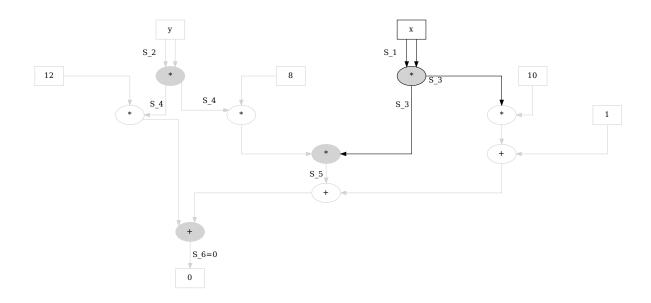
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To compute the number of constraints, observe that we have 3 multiplication gates that have labels on their outgoing edges and 1 addition gate that has a label on its outgoing edge. We therefore have to compute 4 quadratic constraints.

In order to derive the associated R1CS, we have start with an empty R1CS and then iterate over the set $\{S_1, S_2, S_3, S_4, S_5, S_6 = 0\}$ of all edge labels, in order to generate the constraints.

Considering edge label S_1 , we see that the associated edges are not outgoing edges of any algebraic gate, and we therefore have to add no new constraint to the system. The same holds true for edge label S_2 . Looking at edge label S_3 , we see that the associated edges are outgoing edges of a multiplication gate and that the associated subgraph is given by:

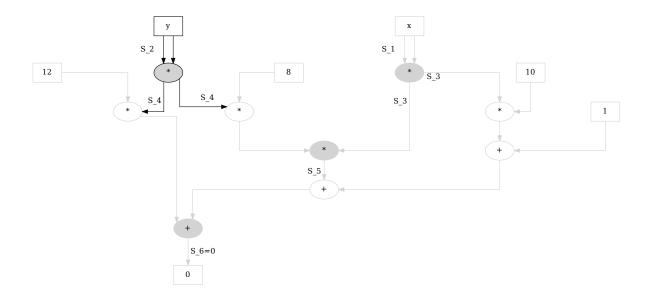


 Both the left and the right input to this multiplication gate are labeled by S_1 . We therefore have to add the following constraint to the system:

$$S_1 \cdot S_1 = S_3$$

Looking at edge label S_4 , we see that the associated edges are outgoing edges of a multiplication gate and that the associated subgraph is given by:

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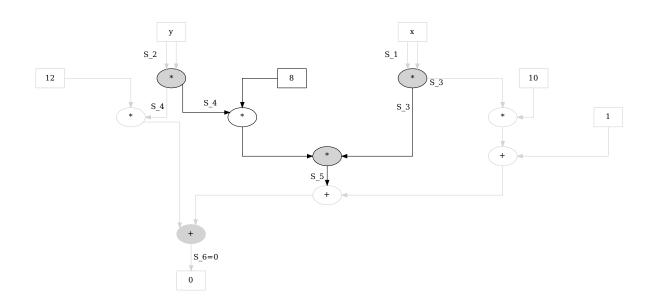
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Both the left and the right input to this multiplication gate are labeled by S_2 and we therefore have to add the following constraint to the system:

$$S_2 \cdot S_2 = S_4$$

Edge label S_5 is more interesting. To see if it implies a constraint, we have to construct the associated subgraph first, which consists of all edges and all nodes in all path starting either at a constant input or a labeled edge. We get

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The right input to the associated multiplication gate is given by the labeled edge S_3 . However, the left input is not a labeled edge, but has a labeled edge in one of its path. This implies that we have to add a constraint to the system. To compute the left factor of that constraint, we have to compute the output of subgraph associated to the left edge, which is $8 \cdot W_2$. This gives the constraint

$$(S_4 \cdot 8) \cdot S_3 = S_5$$

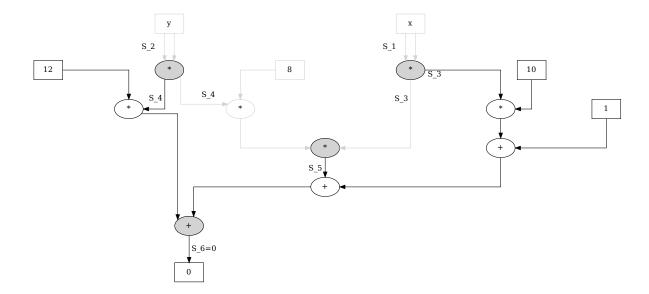
The last edge label is the constant $S_6 = 0$. To see if it implies a constraint, we have to construct the associated subgraph, which consists of all edges and all nodes in all path starting either at a constant input or a labeled edge. We get

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Both the left and the right input are unlabeled, but have a labeled edges in their path. This implies that we have to add a constraint to the system. Since the gate is an addition gate, the right factor in the quadratic constraint is always 1 and the left factor is computed by symbolically executing all inputs to all gates in sub-circuit. We get

$$(12 \cdot S_4 + S_5 + 10 \cdot S_3 + 1) \cdot 1 = 0$$

Since there are no more labeled outgoing edges, we are done deriving the constraints. Combining all constraints together, we get the following R1CS:

$$S_1 \cdot S_1 = S_3$$

$$S_2 \cdot S_2 = S_4$$

$$(S_4 \cdot 8) \cdot S_3 = S_5$$

$$(12 \cdot S_4 + S_5 + 10 \cdot S_3 + 1) \cdot 1 = 0$$

which is equivalent to the R1CS we derived in example 114. The languages $L_{3.fac_zk}$ and $L_{3.fac_circ}$ are therefore equivalent and both the circuit as well as the R1CS are just two different ways to express the same language.

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6.2.3 Quadratic Arithmetic Programs

We have introduced algebraic circuits and their associated rank-1 constraints systems as two particular models able to represent space- and time-bounded computation. Both models define formal languages, and associated membership as well as knowledge claims can be constructively proofed by executing the circuit in order to compute solutions to its associated R1CS.

One reason why those systems are useful in the context of succinct zero-knowledge proof systems is because any R1CS can be transformed into another computational model called **quadratic arithmetic programs** (QAP), which serve as the basis for some of the most efficient succinct non-interactive zero-knowledge proof generators that currently exist.

As we will see, proving statements for languages that have checking relations defined by quadratic arithmetic programs can be achieved by providing certain polynomials, and those proofs can be verified by checking a particular divisibility property.

QAP representation To understand what quadratic arithmetic programs are in detail, let \mathbb{F} be a field and R a rank-1 constraints system over \mathbb{F} such that the number of non-zero elements in \mathbb{F} is strictly larger then the number k of constraints in R. Moreover, let a_j^i , b_j^i and $c_j^i \in \mathbb{F}$ for every index $0 \le j \le n+m$ and $1 \le i \le k$, be the defining constants of the R1CS and m_1, \ldots, m_k be arbitrary, invertible and distinct elements from \mathbb{F} .

Then a **quadratic arithmetic program** [QAP] of the R1CS is the following set of polynomials over \mathbb{F} :

$$QAP(R) = \left\{ T \in \mathbb{F}[x], \left\{ A_j, B_j, C_j \in \mathbb{F}[x] \right\}_{h=0}^{n+m} \right\}$$
 (6.12)

In the equation above, $T(x) := \prod_{l=1}^{k} (x - m_l)$ is a polynomial of degree k, called the **target polynomial** of the QAP and A_j , B_j as well as C_j are the unique degree k-1 polynomials defined by the following equation:

$$A_j(m_i) = a_j^i \quad B_j(m_i) = b_j^i \quad C_j(m_i) = C_j^i \quad j = 1, \dots, n + m + 1, i = 1, \dots, k$$
 (6.13)

Given some rank-1 constraint system, an associated quadratic arithmetic program is therefore nothing but a set of polynomials, computed from the constants in the R1CS. To see that the polynomials A_j , B_j and C_j are uniquely defined by the equations in XXX, recall that a a polynomial of degree k-1 is completely determined on k evaluation points and the equation 4 precisely determines those k evaluation points.

Since we only consider polynomials over fields, Lagrange's interpolation method from 3.30 in chapter 3 can be used to derive the polynomials A_j , B_j and C_j from their defining equations XXX. A practical method to compute a QAP from a given R1CS therefore consists of two steps. If the R1CS consists of k constraints, first choose k invertible and mutually different points from the underlaying field. Every choice defines a different QAP for the same R1CS. Then use Lagrange's method and equation XXX to compute the polynomials A_j , B_j and C_j for every $1 \le j \le k$.

Example 122 (Generalized factorization SNARK). To provide a better intuition of quadratic arithmetic programs and how they are computed from their associated rank-1 constraint systems, consider the language $L_{3.fac.zk}$ from example 106 and its associated R1CS:

$$W_1 \cdot W_2 = W_4$$
 constraint 1
 $W_4 \cdot W_3 = I_1$ constraint 2

In this example we want to transform this R1CS into an associated QAP. In a first step, we have to compute the defining constants a_i^i , b_i^i and c_i^i of the R1CS. According to XXX, we have

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Since the R1CS is defined over the field \mathbb{F}_{13} and has two constraining equations, we need to choose two arbitrary but distinct elements m_1 and m_2 from \mathbb{F}_{13} . We choose $m_1 = 5$, and $m_2 = 7$ and with this choice we get the target polynomial

$$T(x) = (x - m_1)(x - m_2)$$
 # Definition of T
 $= (x - 5)(x - 7)$ # Insert our choice
 $= (x + 8)(x + 6)$ # Negatives in \mathbb{F}_{13}
 $= x^2 + x + 9$ # expand

Then we have to compute the polynomials A_j , B_j and C_j by their defining equation from the R1CS coefficients. Since the R1CS has two constraining equations, those polynomials are of degree 1 and they are defined by their evaluation at the point $m_1 = 5$ and the point $m_2 = 7$.

At point m_1 , each polynomial A_j is defined to be a_j^1 and at point m_2 , each polynomial A_j is defined to be a_j^2 . The same holds true for the polynomials B_j as well as C_j . Writing all these equations now, we get:

$$A_0(5) = 0$$
, $A_1(5) = 0$, $A_2(5) = 1$, $A_3(5) = 0$, $A_4(5) = 0$, $A_5(5) = 0$
 $A_0(7) = 0$, $A_1(7) = 0$, $A_2(7) = 0$, $A_3(7) = 0$, $A_4(7) = 0$, $A_5(7) = 1$
 $B_0(5) = 0$, $B_1(5) = 0$, $B_2(5) = 0$, $B_3(5) = 1$, $B_4(5) = 0$, $B_5(5) = 0$
 $B_0(7) = 0$, $B_1(7) = 0$, $B_2(7) = 0$, $B_3(7) = 0$, $B_4(7) = 1$, $B_5(7) = 0$
 $C_0(5) = 0$, $C_1(5) = 0$, $C_2(5) = 0$, $C_3(5) = 0$, $C_4(5) = 0$, $C_5(5) = 1$
 $C_0(7) = 0$, $C_1(7) = 1$, $C_2(7) = 0$, $C_3(7) = 0$, $C_4(7) = 0$, $C_5(7) = 0$

Lagrange's interpolation implies that a polynomial of degree k, that is, that zero on k+1 points has to be the zero polynomial. Since our polynomials are of degree 1 and determined on 2 points, we therefore know that the only non-zero polynomials in our QAP are A_2 , A_5 , B_3 , B_4 , C_1 and C_5 , and that we can use Lagrange's interpolation to compute them.

To compute A_2 we note that the set S in our version of Lagrange's method is given by $S = \{(x_0, y_0), (x_1, y_1)\} = \{(5, 1), (7, 0)\}$. Using this set we get:

$$A_{2}(x) = y_{0} \cdot l_{0} + y_{1} \cdot l_{1}$$

$$= y_{0} \cdot \left(\frac{x - x_{1}}{x_{0} - x_{1}}\right) + y_{1} \cdot \left(\frac{x - x_{0}}{x_{1} - x_{0}}\right) = 1 \cdot \left(\frac{x - 7}{5 - 7}\right) + 0 \cdot \left(\frac{x - 5}{7 - 5}\right)$$

$$= \frac{x - 7}{-2} = \frac{x - 7}{11}$$

$$= 6(x - 7) = 6x + 10$$
11⁻¹ = 6
- 7 = 6 and 6 \cdot 6 = 10

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To compute A_5 , we note that the set S in our version of Lagrange's method is given by $S = \{(x_0, y_0), (x_1, y_1)\} = \{(5, 0), (7, 1)\}$. Using this set we get:

$$A_5(x) = y_0 \cdot l_0 + y_1 \cdot l_1$$

$$= y_0 \cdot \left(\frac{x - x_1}{x_0 - x_1}\right) + y_1 \cdot \left(\frac{x - x_0}{x_1 - x_0}\right) = 0 \cdot \left(\frac{x - 7}{5 - 7}\right) + 1 \cdot \left(\frac{x - 5}{7 - 5}\right)$$

$$= \frac{x - 5}{2}$$

$$= 7(x - 5) = 7x + 4$$
$-5 = 8$ and $7 \cdot 8 = 4$

Using Lagrange's interpolation, we can deduce that $A_2 = B_3 = C_5$ as well as $A_5 = B_4 = C_1$, since they are polynomials of degree 1 that evaluate to same values on 2 points. Using this, we get the following set of polynomials

$A_0(x) = 0$	$B_0(x) = 0$	$C_0(x) = 0$
$A_1(x) = 0$	$B_1(x) = 0$	$C_1(x) = 7x + 4$
$A_2(x) = 6x + 10$	$B_2(x)=0$	$C_2(x)=0$
$A_3(x) = 0$	$B_3(x) = 6x + 10$	$C_3(x)=0$
$A_4(x) = 0$	$B_4(x) = 7x + 4$	$C_4(x)=0$
$A_5(x) = 7x + 4$	$B_5(x)=0$	$C_5(x) = 6x + 10$

We can invoke Sage to verify our computation. In sage every polynomial ring has a function lagrange_polynomial that takes the defining points as inputs and the associated Lagrange polynomial as output.

Combining this computation with the target polynomial we derived earlier, a quadratic arithmetic program associated to the rank-1 constraint system $R_{3.fac_zk}$ is given as follows:

$$QAP(R_{3.fac_zk}) = \{x^2 + x + 9, \{0,0,6x+10,0,0,6x+10,7x+4,0\}, \{0,7x+4,0,0,0,6x+10\}\}$$
(6.14)

QAP Satisfiability One of the major points of quadratic arithmetic programs in proofing systems is that solutions of their associated rank-1 constraints systems are in 1:1 correspondence with certain polynomials *P* such that *P* is divisible by the target polynomial *T* of the QAP if and only if the solution id a solution. Verifying solutions to the R1CS and hence, checking proper circuit execution is then achievable by polynomial division of *P* by *T*.

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To be more specific, let R be some rank-1 constraints system with associated assignment variables $(I_1, \ldots, I_n; W_1, \ldots, W_m)$ and let QAP(R) be a quadratic arithmetic program of R. Then

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the tuple $(I_1, ..., I_n; W_1, ..., W_m)$ is a solution to the R1CS if and only if the following polynomial is divisible by the target polynomial T:

$$P_{(I;W)} = (A_0 + \sum_{j=1}^{n} I_j \cdot A_j + \sum_{j=1}^{m} W_j \cdot A_{n+j}) \cdot (B_0 + \sum_{j=1}^{n} I_j \cdot B_j + \sum_{j=1}^{m} W_j \cdot B_{n+j}) - (C_0 + \sum_{j=1}^{n} I_j \cdot C_j + \sum_{j=1}^{m} W_j \cdot C_{n+j})$$
(6.15)

Every tuple (I; W) defines a polynomial $P_{(I;W)}$, and, since each polynomial A_j , B_j and C_j is of degree k-1, $P_{(I:W)}$ is of degree $(k-1) \cdot (k-1) = k^2 - 2k + 1$.

To understand how quadratic arithmetic programs define formal languages, observe that every QAP over a field \mathbb{F} defines a decision function over the alphabet $\Sigma_I \times \Sigma_W = \mathbb{F} \times \mathbb{F}$ in the following way:

$$R_{QAP}: \mathbb{F}^* \times \mathbb{F}^* \to \{true, false\}; (I; W) \mapsto \begin{cases} true & P_{(I;W)} \text{ is divisible by } T \\ false & else \end{cases}$$
 (6.16)

This means that every QAP defines a formal language, and, if the QAP is associated to an R1CS, it can be shown that the two languages are equivalent. A **statement** is a membership claim "There is a word (I;W) in L_{QAP} ". A proof to this claim is therefore a polynomial $P_{(I;W)}$, which is verified by dividing $P_{(I;W)}$ by T.

Note the structural similarity to the definition of an R1CS in 6.2.1.1 and the different ways of computing proofs in both systems. For circuits and their associated rank-1 constraints systems, a constructive proof consists of a valid assignment of field elements to the edges of the circuit, or the variables in the R1CS. However, in the case of QAPs, a valid proof consists of a polynomial $P_{(I:W)}$.

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To compute a proof for a statement in L_{QAP} given some instance I, a proofer first needs to compute a constructive proof W, e.g. by executing the circuit. With (I;W) at hand, the proofer can then compute the polynomial $P_{(I;W)}$ and publish it as proof.

Verifying a constructive proof in the case of a circuit is achieved by executing the circuit, comparing the result to the given proof, and verifying the same proof in the R1CS picture means checking if the elements of the proof satisfy all equation.

In contrast, verifying a proof in the case of a QAP is done by polynomial division of the proof P by the target polynomial T of the QAP. The proof checks out if and only if P is divisible by T.

Example 123. Consider the quadratic arithmetic program $QAP(R_{3.fac_zk})$ from example XXX, and its associated R1CS from equation 6.1.0.1. To give an intuition of how proofs in the language $L_{QAP(R_{3.fac_zk})}$ lets consider the instance $I_1 = 11$. As we know from example XXX, $(W_1, W_2, W_3, W_5) = (2, 3, 4, 6)$ is a proper witness, since $(I_1; W_1, W_2, W_3, W_5) = (11; 2, 3, 4, 6)$ is a valid circuit assignment and hence, a solution to $R_{3.fac_zk}$ and a constructive proof for language $L_{R_{3.fac_zk}}$.

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In order to transform this constructive proof into a membership proof in language $L_{QAP(R_{3.fac_zk})}$ a proofer has to use the elements of the constructive proof, to compute the polynomial $P_{(I:W)}$.

In the case of $(I_1; W_1, W_2, W_3, W_5) = (11; 2, 3, 4, 6)$, the associated proof is computed as fol-

lows:

$$\begin{split} P_{(I;W)} = & (A_0 + \sum_{j=1}^{n} I_{j} \cdot A_{j} + \sum_{j=1}^{m} W_{j} \cdot A_{n+j}) \cdot (B_0 + \sum_{j=1}^{n} I_{j} \cdot B_{j} + \sum_{j=1}^{m} W_{j} \cdot B_{n+j}) - (C_0 + \sum_{j=1}^{n} I_{j} \cdot C_{j} + \sum_{j=1}^{m} W_{j} \cdot C_{n+j}) \\ = & (2(6x+10) + 6(7x+4)) \cdot (3(6x+10) + 4(7x+4)) - (11(7x+4) + 6(6x+10)) \\ = & ((12x+7) + (3x+11)) \cdot ((5x+4) + (2x+3)) - ((12x+5) + (10x+8)) \\ = & (2x+5) \cdot (7x+7) - (9x) \\ = & (x^2 + 2 \cdot 7x + 5 \cdot 7x + 5 \cdot 7) - (9x) \\ = & (x^2 + x + 9x + 9) - (9x) \\ = & x^2 + x + 9 \end{split}$$

Given instance $I_1 = 11$ a proofer therefore provides the polynomial $x^2 + x + 9$ as proof. To verify this proof, any verifier can then look up the target polynomial T from the QAP and divide $P_{(I;W)}$ by T. In this particular example, $P_{(I;W)}$ is equal to the target polynomial T, and hence, it is divisible by T with P/T = 1. The verification therefore checks the proof.

To give an example of a false proof, consider the tuple $(I_1; W_1, W_2, W_3, W_4) = (11, 2, 3, 4, 8)$. Executing the circuit, we can see that this is not a valid assignment and not a solution to the R1CS, and hence, not a constructive knowledge proof in $L_{3.fac_zk}$. However, a proofer might use these values to construct a false proof $P_{(I:W)}$:

$$P'_{(I;W)} = (A_0 + \sum_{j=1}^{n} I_j \cdot A_j + \sum_{j=1}^{m} W_j \cdot A_{n+j}) \cdot (B_0 + \sum_{j=1}^{n} I_j \cdot B_j + \sum_{j=1}^{m} W_j \cdot B_{n+j}) - (C_0 + \sum_{j=1}^{n} I_j \cdot C_j + \sum_{j=1}^{m} W_j \cdot C_{n+j})$$

$$= (2(6x+10) + 8(7x+4)) \cdot (3(6x+10) + 4(7x+4)) - (8(6x+10) + 11(7x+4))$$

$$= 8x^2 + 6$$

Given instance $I_1 = 11$, a proofer therefore provides the polynomial $8x^2 + 6$ as proof. To verify this proof, any verifier can look up the target polynomial T from the QAP and divide $P_{(I;W)}$ by T. However, polynomial division has the following remainder:

$$(8x^2+6)/(x^2+x+9) = 8 + \frac{5x+12}{x^2+x+9}$$

This implies that $P_{(I;W)}$ is not divisible by T, and hence, the verification does not check the proof. Any verifier can therefore show that the proof is false.

4798	sage: $P == F13t(8*t^2 + 6)$	651
4799	True	652
4800	<pre>sage: P % T # remainder</pre>	653
4801	5*t + 12	654

Chapter 7

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Circuit Compilers

As we have seen in the previous chapter, statements can be formalized as membership or knowledge claims in formal language, and algebraic circuits as well as rank-1 constraint systems are two practically important ways to define those languages.

However, both algebraic circuits and rank-1 constraint systems are not ideal from a developers point of view, because they deviate substantially from common programing paradigms. Writing real-world applications as circuits and the associated verification in terms of rank-1 constraint systems is as least as troublesome as writing any other low-level language like assembler code. To allow for complex statement design, it is therefore necessary to have some kind of compiler framework, capable of transforming high-level languages into arithmetic circuits and associated rank-1 constraint systems.

As we have seen in chapter 6 and in 6.2.1.1,, both arithmetic circuits and rank-1 constraint systems have a modularity property by which it is possible to synthesize complex circuits from simple ones. A basic approach taken by many circuit/R1CS compilers is therefore to provide a library of atomic and simple circuits and then define a way to combine those basic building blocks into arbitrary complex systems.

In this chapter, we provide an introduction to basic concepts of so-called **circuit compilers** and derive a toy language which we can "compile" in a pen-and-paper approach into algebraic circuits and their associated rank-1 constraint systems.

We start with a general introduction to our language, and then introduce atomic types like booleans and unsigned integers. Then we define the fundamental control flow primitives like the if-then-else conditional and the bounded loop. We will look at basic functionality primitives like elliptic curve cryptography. Primitives like these are often called **gadgets** in the literature.

A Pen-and-Paper Language 7.1 4826

To explain basic concepts of circuit compilers and their associated high-level languages, we derive an informal toy language and associated "brain-compiler" which we name PAPER (Pen-4828 And-Paper Execution Rules). PAPER allows programmers to define statements in Rust-like 4829 pseudo-code. The language is inspired by ZOKRATES and circom. 4830

The Grammar 7.1.1

In PAPER, any statement is defined as an ordered list of functions, where any function has to be declared in the list before it is called in another function of that list. The last entry in a statement has to be a special function, called main. Functions take a list of typed parameters as inputs

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add references to these languages?

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and compute a tuple of typed variables as output, where types are special functions that define how to transform that type into another type, ultimately transforming any type into elements of the base field where the circuit is defined over.

Any statement is parameterized over the field that the circuit will be defined on, and has additional optional parameters of unsigned type, needed to define the size of array or the counter of bounded loops. The following definition makes the grammar of a statement precise using a command line language like description:

```
statement <Name> {F:<Field> [ , <N_1: unsigned>,... ] } {
4842
      [fn <Name>([[pub]<Arg>:<Type>,...]) -> (<Type>,...) {
4843
        [let [pub] <Var>:<Type> ;... ]
4844
        [let const <Const>:<Type>=<Value> ; ... ]
4845
        Var<==(fn([<Arg>|<Const>|<Var>,...])|(<Arg>|<Const>|<Var>));
4846
        return (<Var>,...);
      } ; . . . ]
4848
      fn main([[pub]<Arg>:<Type>,...]) -> (<Type>,...){
4849
        [let [pub] <Var>:<Type> ;... ]
4850
        [let const <Const>:<Type>=<Value> ;... ]
4851
        Var<==(fn([<Arg>|<Const>|<Var>,...])|(<Arg>|<Const>|<Var>));
        return (<Var>,...);
4853
      } ;
4854
4855
```

Function arguments and variables are private by default, but can be declared as public by the pub specifier. Declaring arguments and variables as public always overwrites any previous or conflicting private declarations. Every argument, constant or variable has a type, and every type is defined as a function that transforms that type into another type:

```
4860 type <TYPE>( t1 : <TYPE_1>) -> TYPE_2 {
4861    let t2: TYPE_2 <== fn(TYPE_1)
4862    return t2
4863 }
```

Many real-world circuit languages are based on a similar, but of course more sophisticated approach than PAPER. The purpose of PAPER is to show basic principles of circuit compilers and their associated high-level languages.

Example 124. To get a better understanding of the grammar of PAPER, the following constitutes proper high-level code that follows the grammar of the PAPER language, assuming that all types in that code have been defined elsewhere.

```
statement MOCK_CODE \{F: F_43, N_1 = 1024, N_2 = 8\} {
4870
      fn foo(in 1 : F, pub in 2 : TYPE 2) \rightarrow F {
4871
        let const c_1 : F = 0;
4872
        let const c_2 : TYPE_2 = SOME_VALUE ;
4873
        let pub out_1 : F ;
4874
        out 1<== c 1 ;
4875
        return out 1 ;
4876
4877
4878
      fn bar(pub in_1 : F) -> F {
4879
        let out_1 : F ;
4880
        out_1<==foo(in_1);
4881
        return out_1 ;
4882
```

```
} ;
4883
4884
      fn main(in_1 : TYPE_1) -> (F, TYPE_2) {
4885
         let const c_1 : TYPE_1
                                    = SOME_VALUE ;
4886
         let const c_2 : F = 2;
4887
         let const c 3 : TYPE 2
                                     = SOME VALUE
         let pub out_1 : F ;
4889
         let out_2 : TYPE_2 ;
4890
         c_1 <== in_1 ;
4891
         out 1 <== foo(c 2);
4892
         out 2 <== TYPE 2 ;
4893
         return (out_1,out_2) ;
4894
4895
    }
4896
```

7.1.2 The Execution Phases

In contrast to normal executable programs, programs for circuit compilers have two modes of execution. The first mode, usually called the **setup phase**, is executed in order to generate the circuit and its associated rank-1 constraint system, the latter of which is then usually used as input to some zero-knowledge proof system.

The second mode of execution is usually called the **prover phase**. In this phase, a prover usually computes a valid assignment to the circuit. Depending on the use case, this valid assignment is then either directly used as constructive proof for proper circuit execution or is transferred as input to the proof generation algorithm of some zero-knowledge proof system, where the full-sized, non hiding constructive proof is processed into a succinct proof with various levels of zero knowledge.

Modern circuit languages and their associated compilers abstract over those two phases and provide a unified **interphase** to the developer, who then writes a single program that can be used in both phases.

To give the reader a clear, conceptual distinction between the two phases, PAPER keeps them separated. Code can be "brain-compiled" during the **setup-phase** in a pen-and-paper approach into visual circuits. Once a circuit is derived, it can be executed in a **prover phase** to generate a valid assignment. The valid assignment is then interpreted as a constructive proof for a knowledge claim in the associated language.

The Setup Phase In PAPER, the task of the setup phase is to compile code in the PAPER language into a visual representation of an algebraic circuit. Deriving the circuit from the code in a pen-and-paper style is what we call **brain-compiling**.

Given some statement description that adheres to the correct grammar, we start circuit development with an empty circuit, compile the main function first and then inductively compile all other functions as they are called during the process.

For every function we compile, we draw a box-node for every argument, every variable and every constant of that function. If the node represents a variable, we label it with that variable's name, and if it represents a constant, we label it with that constant's value. We group arguments into a subgraph labeled "inputs" and return values into a subgraph labeled "outputs". We then group everything into a subgraph and label that subgraph with the function's name.

After this is done, we have to do a consistency and type check for every occurrence of the

assignment operator <==. We have to ensure that the expression on the right side of the operator is well defined and that the types of both side match.

Then we compile the right side of every occurrence of the assignment operator <==. If the right side is a constant or variable defined in this function, we draw a dotted line from the box-node that represents the left side of <== to the box node that represents the right side of the same operator. If the right side represents an argument of that function we draw a line from the box-node that represents the left side of <== to the box node that represents the right side of the same operator.

If the right side of the <== operator is a function, we look into our database, find its associated circuit and draw it. If no circuit is associated to that function yet, we repeat the compilation process for that function, drawing edges from the function's argument to its input nodes and from the functions output nodes to the nodes on the right side of <==.

During that process, edge labels are drawn according to the rules from 6.2.2.1. If the associated variable represents a private value, we use the W label to indicate a witness, and if it represents a public value, we use the I label to indicate an instance.

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Once this is done, we compile all occurring types in a function, by compiling the function of each type. We do this inductively until we reach the type of the base field. Circuits have no notion of types, only of field elements; hence, every type needs to be compiled to the field type in a sequence of compilation steps.

The compilation stops once we have inductively replaced all functions by their circuits. The result is a circuit that contains many unnecessary box nodes. In a final optimization step, all box nodes that are directly linked to each other are collapsed into a single node, and all box nodes that represent the same constants are collapsed into a single node.

Of course, PAPER's brain-compiler is not properly defined in any formal manner. Its purpose is to highlight important steps that real-world compilers undergo in their setup phases.

Example 125 (A trivial Circuit). To give an intuition of how to write and compile circuits in the PAPER language, consider the following statement description:

```
statement trivial circuit {F:F 13}
4955
      fn main\{F\} (in1 : F, pub in2 : F) -> (F,F)\{F\}
4956
         let const outc1 : F = 0;
4957
         let const inc1 : F = 7;
4958
         let out1 : F ;
4959
         let out2 : F ;
4960
         out1 <== inc1;
4961
         out2 <== in1;
4962
         outc1 <== in2;
4963
         return (out1, out2) ;
      }
4965
    }
4966
```

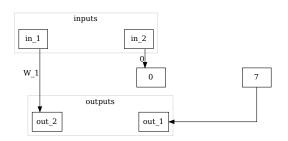
To brain-compile this statement into an algebraic circuit with PAPER, we start with an empty circuit and evaluate function main, which is the only function in this statement.

We draw box-nodes for every argument, every constant and every variable of the function and label them with their names or values, respectively. Then we do a consistency and type check for every <== operator in the function. Since the circuit only wires inputs to outputs and all elements have the same type, the check is valid.

Then we evaluate the right side of the assignment operators. Since, in our case, the right side of each operator is not a function, we draw edges from the box-nodes on the right side to the associated box node on the left side. To label those edges, we use the general rules of

algebraic circuits as defined in 6.2.2.1. According to those rules, every incoming edge of a sink node has a label and every outgoing edge of a source node has a label, if the node is labeled with a variable. Since nodes that represent constants are implicitly assumed to be private, and since the public specifier determines if an edge is labeled with *W* or *I*, we get the following circuit:

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The Prover Phase In PAPER, a so-called prover phase can be executed once the setup phase has generated a circuit image from its associated high-level code. This is done by executing the circuit while assigning proper values to all input nodes of the circuit. However, in contrast to most real-world compilers, PAPER does not tell the prover how to find proper input values to a given circuit. Real-world programing languages usually provide this data by computations that are done outside of the circuit.

Example 126. Consider the circuit from example 125. Valid assignments to this circuit are constructive proofs that the pair of inputs (S_1, S_2) is a point on the tiny-jubjub curve. However, the circuit does not provide a way to actually compute proper values for S_1 and S_2 . Any real-world system therefore needs an auxiliary computation that provides those values.

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7.2 Common Programing concepts

In this section, we cover concepts that appear in almost every programming language, and see how they can be implemented in circuit compilers.

7.2.1 Primitive Types

Primitive data types like booleans, (unsigned) integers, or strings are the most basic building blocks one can expect to find in every general high-level programing language. In order to write statements as computer programs that compile into circuits, it is therefore necessary to implement primitive types as constraint systems, and define their associated operations as circuits.

In this section, we look at some common ways to achieve this. After a recapitulation of the atomic type of prime field elements, we start with an implementation of the boolean type and its associated boolean algebra as circuits. After that, we define unsigned integers based on the boolean type, and leave the implementation of signed integers as an exercise to the reader.

It should be noted, however, that while primitive data types in common programing languages (like C, Go, or Rust) have a one-to-one correspondence with objects in the computer's memory, this is not the case for most languages that compile into algebraic circuits. As we will see in the following paragraphs, common primitives like booleans or unsigned integers require many constraints and memory. Primitives different from the underlying field elements can be expensive.

5009 The base-field type

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Since both algebraic circuits and their associated rank-1 constraint systems are defined over a finite field, elements from that field are the atomic informational units in those models. In this sense, field elements $x \in \mathbb{F}$ are for algebraic circuits what bits are for computers.

In PAPER, we write F for this type and specify the actual field instance for every statement in curly brackets after the name of that statement. Two functions are associated to this type, which are induced by the **addition** and **multiplication** law in the field F. We write

$$MUL: F \times F \to F; (x,y) \mapsto MUL(x,y)$$
(7.1)

ADD: $F \times F \to F$; $(x, y) \mapsto ADD(x, y)$ (7.2)

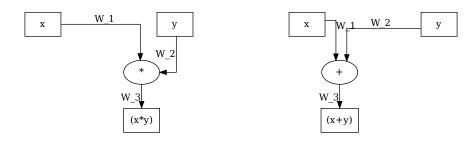
Circuit compilers have to compile these functions into algebraic gates, as explained in 6.2.2.2. Every other function has to be expressed in terms of them and proper wiring.

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To represent addition and multiplication in the PAPER language, we define the following two functions:

```
5021 fn MUL(x : F, y : F) \rightarrow (MUL(x,y):F) {}
5022 fn ADD(x : F, y : F) \rightarrow (ADD(x,y):F) {}
```

The compiler then compiles every occurrence of the MUL or the ADD function into the following circuits:



Example 127 (Basic gates). To give an intuition of how a real-world compiler might transform addition and multiplication in algebraic expressions into a circuit, consider the following PAPER statement:

```
5029  statement basic_ops {F:F_13} {
5030    fn main(in_1 : F, pub in_2 : F) -> (out_1:F, out_2:F) {
5031       out_1 <== MUL(in_1,in_2) ;
5032       out_2 <== ADD(in_1,in_2) ;
5033    }
5034 }</pre>
```

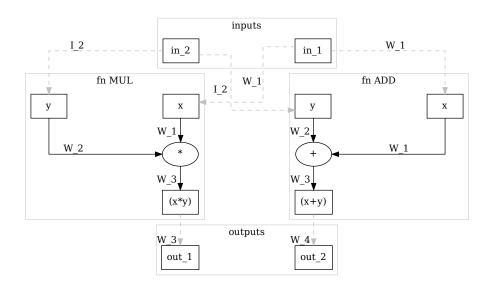
To compile this into an algebraic circuit, we start with an empty circuit and evaluate the function main, which is the only function in this statement.

We draw an inputs subgraph containing box-nodes for every argument of the function, and an outputs subgraph containing box-nodes for every factor in the return value. Since all of these nodes represent variables of the field type, we don't have to add any type constraints to the circuit.

We check the validity of every expression on the right side of every <== operator including a type check. In our case, every variable is of the field type and hence the types match the types of the MUL as well as the ADD function and the type of the left sides of <== operators.

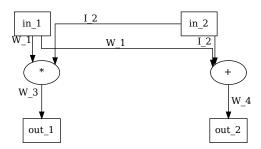
We evaluate the expressions on the right side of every <== operator inductively, replacing every occurrence of a function with a subgraph that represents its associated circuit.

According to PAPER, every occurrence of the public specifier overwrites the associate private default value. Using the appropriate edge labels we get:



Any real-world compiler might process its associated high-level language in a similar way, replacing functions, or gadgets by predefined associated circuits. This process is often followed by various optimization steps that try to reduce the number of constraints as much as possible.

In PAPER, we optimize this circuit by collapsing all box nodes that are directly connected to other box nodes, adhering to the rule that a variable's public specifier overwrites any private specifier. Reindexing edge labels, we get the following circuit as our pen and pencil compiler output:



Example 128 (3-factorization). Consider our 3-factorization problem from example 106_and the associated circuit $C_{3.fac_zk}(\mathbb{F}_{13})$ we provided in equation 6.8. To understand the process of replacing high-level functions by their associated circuits inductively, we want define a PAPER statement that we brain-compile into an algebraic circuit equivalent to $C_{3.fac_zk}(\mathbb{F}_{13})$:

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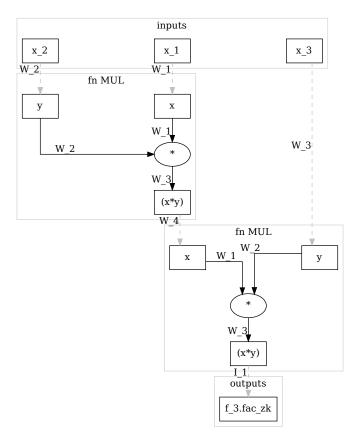
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```
5061 statement 3_fac_zk {F:F_13} {
5062    fn main(x_1 : F, x_2 : F, x_3 : F) -> (pub 3_fac_zk : F) {
5063       f_3.fac_zk <== MUL( MUL( x_1 , x_2 ) , x_3 ) ;
5064    }
5065 }</pre>
```

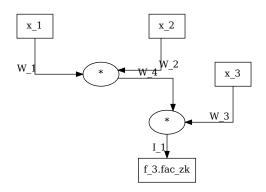
Using PAPER, we start with an empty circuit and then add 3 input nodes to the input subgraph as well as 1 output node to the output subgraph. All these nodes are decorated with the associated variable names. Since all of these nodes represent variables of the field type, we don't have to add any type constraints to the circuit.

We check the validity of every expression on the right side of the single <== operator including a type check.

We evaluate the expressions on the right side of every <== operator inductively. We have two nested multiplication functions and we replace them by the associated multiplication circuits, starting with the most outer function. We get:

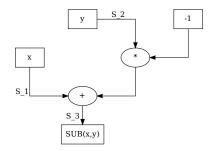


In a final optimization step, we collaps all box nodes directly connected to other box nodes, adhering to the rule that a variables public specifier overwrites any private specifier. Reindexing edge labels we get the following circuit:



The Subtraction Constraint System By definition, algebraic circuits only contain addition and multiplication gates, and it follows that there is no single gate for field subtraction, despite the fact that subtraction is a native operation in every field.

High-level languages and their associated circuit compilers, therefore, need another way to deal with subtraction. To see how this can be achieved, recall that subtraction is defined by addition with the additive inverse, and that the inverse can be computed efficiently by multiplication with -1. A circuit for field subtraction is therefore given by



Using the general method from 6.2.1.1, the circuits associated rank-1 constraint system is given by:

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$$(S_1 + (-1) \cdot S_2) \cdot 1 = S_3 \tag{7.3}$$

Any valid assignment $\{S_1, S_2, S_3\}$ to this circuit therefore enforces the value S_3 to be the difference $S_1 - S_2$.

Real-world compilers usually provide a gadget or a function to abstract over this circuit such that programers can use subtraction as if it were native to circuits. In PAPER, we define the following subtraction function that compiles to the previous circuit:

```
5095 fn SUB(x : F, y : F) -> (SUB(x,y) : F) {
5096    constant c : F = -1;
5097    SUB <== ADD(x , MUL(y , c));
5098  }
```

In the setup phase of a statement, we compile every occurrence of the SUB function into an instance of its associated subtraction circuit, and edge labels are generated according to the rules from 6.2.2.1.

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The Inversion Constraint System By definition, algebraic circuits only contain addition and multiplication gates, and it follows that there is no single gate for field inversion, despite the fact that inversion is a native operation in every field.

If the underlying field is a prime field, one approach would be to use Fermat's little theorem 3.3 to compute the multiplicative inverse inside the circuit. To see how this works, let \mathbb{F}_p be the prime field. The multiplicative inverse x^{-1} of a field element $x \in \mathbb{F}$ with $x \neq 0$ is then given by $x^{-1} = x^{p-2}$, and computing x^{p-2} in the circuit therefore computes the multiplicative inverse.

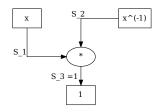
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Unfortunately, real-world primes p are large and computing x^{p-2} by repeated multiplication of x with itself is infeasible. A "double and multiply" approach (as described in XXX) is faster, as it computes the power in roughly $log_2(p)$ steps, but still adds a lot of constraints to the circuit.

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Computing inverses in the circuit makes no use of the fact that inversion is a native operation in any field. A more constraints friendly approach is therefore to compute the multiplicative inverse outside of the circuit and then only enforce correctness of the computation in the circuit.

To understand how this can be achieved, observe that a field element $y \in \mathbb{F}$ is the multiplicative inverse of a field element $x \in \mathbb{F}$ if and only if $x \cdot y = 1$ in \mathbb{F} . We can use this, and define a circuit that has two inputs, x and y, and enforces $x \cdot y = 1$. It is then guaranteed that y is the multiplicative inverse of x. The price we pay is that we can not compute y by circuit execution, but auxiliary data is needed to tell any prover which value of y is needed for a valid circuit assignment. The following circuit defines the constraint



Using the general method from 6.2.1.1, the circuit is transformed into the following rank-1 constraint system:

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$$S_1 \cdot S_2 = 1 \tag{7.4}$$

Any valid assignment $\{S_1, S_2\}$ to this circuit enforces that S_2 is the multiplicative inverse of S_1 , and, since there is no field element S_2 such that $0 \cdot S_2 = 1$, it also handles the fact that the multiplicative inverse of 0 is not defined in any field.

Real-world compilers usually provide a gadget or a function to abstract over this circuit, and those functions compute the inverse x^{-1} as part of their witness generation process. Programers then don't have to care about providing the inverse as auxiliary data to the circuit. In PAPER, we define the following inversion function that compiles to the previous circuit:

```
fn INV(x : F, y : F) -> (x_inv : F) {
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fn INV(x : F, y : F) -> (x_
```

As we see, this functions takes two inputs, the field value and its inverse. It therefore does not handle the computation of the inverse by itself. This is to keep PAPER as simple as possible.

In the setup phase, we compile every occurrence of the INV function into an instance of the inversion circuit XXX, and edge labels are generated according to the rules from XXX.

add reference

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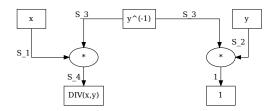
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The Division Constraint System By definition, algebraic circuits only contain addition and multiplication gates, and it follows that there is no single gate for field division, despite the fact that division is a native operation in every field.

Implementing division as a circuit, we use the fact that division is multiplication with the multiplicative inverse. We therefore define division as a circuit using the inversion circuit and constraint system from the previous paragraph. Expensive inversion is computed outside of the circuit and then provided as circuit input. We get



Using the method from 6.2.1.1, we transform this circuit into the following rank-1 constraint system:

check reference

$$S_2 \cdot S_3 = 1$$
$$S_1 \cdot S_3 = S_4$$

Any valid assignment $\{S_1, S_2, S_3, S_4\}$ to this circuit enforces S_4 to be the field division of S_1 by S_2 . It handles the fact that division by 0 is not defined, since there is no valid assignment in case $S_2 = 0$.

In PAPER, we define the following division function that compiles to the previous circuit:

```
5152 fn DIV(x : F, y : F, y_inv : F) -> (DIV : F) {
5153    DIV <== MUL(x, INV(y, y_inv));
5154 }
```

In the setup phase, we compile every occurrence of the binary INV operator into an instance of the inversion circuit.

Exercise 44. Let F be the field \mathbb{F}_5 of modular 5 arithmetics from example 13. Brain-compile the following PAPER statement into an algebraic circuit:

check reference

```
statement STUPID_CIRC {F: F_5} {
5159
      fn foo(in_1 : F, in_2 : F) -> (out_1 : F, out_2 : F,) {
5160
        constant c_1 : F = 3;
5161
        out_1<== ADD( MUL( c_1 , in_1 ) , in_1 ) ;
5162
        out_2<== INV( c_1 , in_2 ) ;
5163
      } ;
5164
5165
      fn main(in_1 : F, in_2 ; F) -> (out_1 : F, out_2 : TYPE_2) {
5166
        constant (c_1, c_2) : (F, F) = (3, 2);
5167
         (out_1, out_2) <== foo(in_1, in_2);
5168
      } ;
5169
    }
5170
```

Exercise 45. Consider the tiny-jubjub curve from example 66 and its associated circuit 125.
Write a statement in PAPER that brain-compiles the statement into a circuit equivalent to the one derived in XXX, assuming that curve points are instances and every other assignment is a witness.

check reference

check reference

Exercise 46. Let $F = \mathbb{F}_{13}$ be the modular 13 prime field and $x \in F$ some field element. Define a statement in PAPER such that given instance x a field element $y \in F$ is a witness for the statement if and only if y is the square root of x.

Brain-compile the statement into a circuit and derive its associated rank-1 constraint system. Consider the instance x = 9 and compute a constructive proof for the statement.

The boolean Type

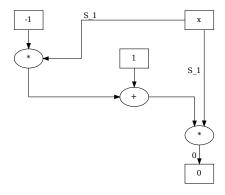
Booleans are a classical primitive type, implemented by virtually every higher programing language. It is therefore important to implement booleans in circuits. One of the most common ways to do this is by interpreting the additive and multiplicative neutral element $\{0,1\} \subset \mathbb{F}$ as the two boolean values such that 0 represents *false* and 1 represents *true*. boolean operators like *and*, *or*, or *xor* are then expressible as algebraic computations inside \mathbb{F} .

Representing booleans this way is convenient, because the elements 0 and 1 are defined in any field. The representation is therefore independent of the actual field in consideration.

To fix boolean algebra notation, we write 0 to represent *false* and 1 to represent *true*, and we write \land to represent the boolean AND as well as \lor to represent the boolean OR operator. The boolean NOT operator is written as \neg .

The boolean Constraint System To represent booleans by the additive and multiplicative neutral elements of a field, a constraint is required to actually enforce variables of boolean type to be either 1 or 0. In fact, many of the following circuits that represent boolean functions are only correct under the assumption that their input variables are constrained to be either 0 or 1. Not constraining boolean variables is a common problem in circuit design.

In order to constrain an arbitrary field element $x \in \mathbb{F}$ to be 1 or 0, the key observation is that the equation $x \cdot (1-x) = 0$ has only two solutions 0 and 1 in any field. Implementing this equation as a circuit therefore generates the correct constraint:



Using the method from 6.2.1.1, we transform this circuit into the following rank-1 constraint system:

check reference

$$S_1 \cdot (1 - S_1) = 0$$

Any valid assignment $\{S_1\}$ to this circuit enforces S_1 to be either 0 or 1.

Some real-world circuit compilers (like ZOKRATES or BELLMAN) are typed, while others (like circom) are not. However, all of them have their way of dealing with the binary constraint. In PAPER, we define the following boolean type that compiles to the previous circuit:

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```
type BOOL(b : BOOL) -> (x : F) {
    constant c1 : F = 0 ;
    constant c2 : F = 1 ;
    constant c3 : F = -1 ;
    constant c4 : F = 0 ;
    constant c5 : F = 0 ;
    c
```

In the setup phase of a statement, we compile every occurrence of a variable of boolean type into an instance of its associated boolean circuit.

The AND operator constraint system Given two field elements b_1 and b_2 from \mathbb{F} that are constrained to represent boolean variables, we want to find a circuit that computes the logical and operator $AND(b_1,b_2)$ as well as its associated R1CS that enforces b_1 , b_2 , $AND(b_1,b_2)$ to satisfy the constraint system if and only if $b_1 \wedge b_2 = AND(b_1,b_2)$ holds true.

The key insight here is that, given three boolean constraint variables b_1 , b_2 and b_3 , the equation $b_1 \cdot b_2 = b_3$ is satisfied in \mathbb{F} if and only if the equation $b_1 \wedge b_2 = b_3$ is satisfied in boolean algebra. The logical operator \wedge is therefore implementable in \mathbb{F} by field multiplication of its arguments and the following circuit computes the \wedge operator in \mathbb{F} , assuming all inputs are restricted to be 0 or 1:



The associated rank-1 constraint system can be deduced from the general process 6.2.1.1 and consists of the following constraint:

check reference

$$S_1 \cdot S_2 = S_3 \tag{7.5}$$

Common circuit languages typically provide a gadget or a function to abstract over this circuit such that programers can use the \land operator without caring about the associated circuit. In PAPER, we define the following function that compiles to the \land -operator's circuit:

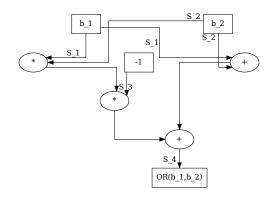
```
5228 fn AND(b_1 : BOOL, b_2 : BOOL) -> AND(b_1,b_2) : BOOL{
5229 AND(b_1,b_2) <== MUL(b_1, b_2);
5230 }
```

In the setup phase of a statement, we compile every occurrence of the AND function into an instance of its associated ∧-operator's circuit.

The OR operator constraint system Given two field elements b_1 and b_2 from \mathbb{F} that are constrained to represent boolean variables, we want to find a circuit that computes the logical **or** operator $OR(b_1, b_2)$ as well as its associated R1CS that enforces b_1 , b_2 , $OR(b_1, b_2)$ to satisfy the constraint system if and only if $b_1 \vee b_2 = OR(b_1, b_2)$ holds true.

Assuming that three variables b_1 , b_2 and b_3 are boolean constraint, the equation $b_1 + b_2 - b_1 \cdot b_2 = b_3$ is satisfied in \mathbb{F} if and only if the equation $b_1 \vee b_2 = b_3$ is satisfied in boolean algebra. The logical operator \vee is therefore implementable in \mathbb{F} by the following circuit, assuming all inputs are restricted to be 0 or 1:

"constraints or "constrained"?



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The associated rank-1 constraint system can be deduced from the general process 6.2.1.1 and consists of the following constraints:

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$$S_1 \cdot S_2 = S_3$$
$$(S_1 + S_2 - S_3) \cdot 1 = S_4$$

Common circuit languages typically provide a gadget or a function to abstract over this circuit such that programers can use the \lor operator without caring about the associated circuit. In PAPER, we define the following function that compiles to the \lor -operator's circuit:

```
5245 fn OR(b_1 : BOOL, b_2 : BOOL) -> OR(b_1,b_2) : BOOL{
5246 constant c1 : F = -1 ;
5247 OR(b_1,b_2) <== ADD(ADD(b_1,b_2), MUL(c1, MUL(b_1,b_2))) ;
5248 }
```

In the setup phase of a statement, we compile every occurrence of the OR function into an instance of its associated \vee -operator's circuit.

Exercise 47. Let $\mathbb F$ be a finite field and let b_1 as well as b_2 two boolean constraint variables from $\mathbb F$. Show that the equation $OR(b_1,b_2)=1-(1-b_1)\cdot(1-b_2)$ holds true.

Use this equation to derive an algebraic circuit with ingoing variables b_1 and b_2 and outgoing variable $OR(b_1, b_2)$ such that b_1 and b_2 are boolean constraint and the circuit has a valid assignment, if and only if $OR(b_1, b_2) = b_1 \vee b_2$.

Use the technique from XXX to transform this circuit into a rank-1 constraint system and find its full solution set. Define a PAPER function that brain-compiles into the circuit.

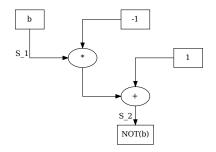
"constraints or "constrained"?

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The NOT operator constraint system Given a field element b from \mathbb{F} that is constrained to represent a boolean variable, we want to find a circuit that computes the logical **NOT** operator NOT(b) as well as its associated R1CS that enforces b, NOT(b) to satisfy the constraint system if and only if $\neg b = NOT(b)$ holds true.

Assuming that two variables b_1 and b_2 are boolean constraint, the equation $(1-b_1) = b_2$ is satisfied in \mathbb{F} if and only if the equation $\neg b_1 = b_2$ is satisfied in boolean algebra. The logical operator \neg is therefore implementable in \mathbb{F} by the following circuit, assuming all inputs are restricted to be 0 or 1:

"constraints or "constrained"?



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The associated rank-1 constraint system can be deduced from the general process XXX_and_consists of the following constraints

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$$(1-S_1)\cdot 1=S_2$$

Common circuit languages typically provide a gadget or a function to abstract over this circuit such that programers can use the ¬ operator without caring about the associated circuit. In PAPER, we define the following function that compiles to the ¬-operator's circuit:

```
5270 fn NOT(b : BOOL -> NOT(b) : BOOL{
5271    constant c1 = 1 ;
5272    constant c2 = -1 ;
5273    NOT(b_1) <== ADD(c1 , MUL(c2 , b));
5274 }</pre>
```

In the setup phase of a statement, we compile every occurrence of the NOT function into an instance of its associated ¬-operator's circuit.

Exercise 48. Let \mathbb{F} be a finite field. Derive the algebraic circuit and associated rank-1 constraint system for the following operators: NOR, XOR, NAND, EQU.

Modularity As we have seen in chapter 6, both algebraic circuits and R1CS have a modularity property, and as we have seen in this section, all basic boolean functions are expressible in circuits. Combining those two properties, show that it is possible to express arbitrary boolean functions as algebraic circuits.

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This shows that the expressiveness of algebraic circuits and therefore rank-1 constraint systems is as general as the expressiveness of boolean circuits. An important implication is that the languages $L_{R1CS-SAT}$ and $L_{Circuit-SAT}$ as defined in 2, are as general as the famous language L_{3-SAT} , which is known to be \mathscr{N} -complete.

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Example 129. To give an example of how a compiler might construct complex boolean expressions in algebraic circuits from simple ones and how we derive their associated rank-1 constraint systems, let's look at the following PAPER statement:

```
5290  statement BOOLEAN_STAT {F: F_p} {
5291    fn main(b_1:BOOL,b_2:BOOL,b_3:BOOL,b_4:BOOL) -> pub b_5:BOOL {
5292      b_5 <== AND(OR(b_1, b_2), AND(b_3, NOT(b_4)));
5293    };
5294 }</pre>
```

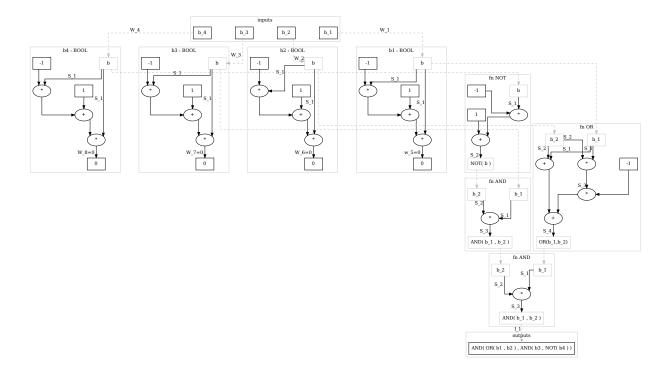
The code describes a circuit that takes four private inputs b_1 , b_2 , b_3 and b_4 of boolean type and computes a public output b_5 such that the following boolean expression holds true:

$$(b_1 \lor b_2) \land (b_3 \land \neg b_4) = b_5$$

During a setup-phase, a circuit compiler transforms this high-level language statement into a circuit and associated rank-1 constraint systems and hence defines a language $L_{BOOLEAN\ STAT}$.

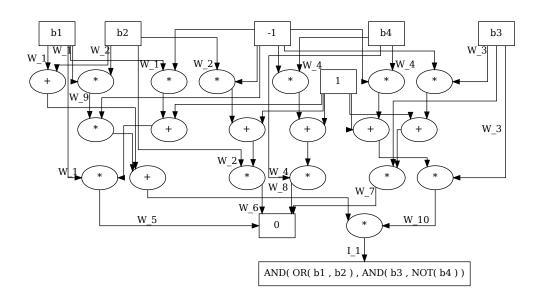
To see how this might be achieved, we use PAPER as an example to execute the setup-phase and compile BOOLEAN_STAT into a circuit. Taking the definition of the boolean constraint XXX as well as the definitions of the appropriate boolean operators into account, we get the following circuit:

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Simple optimization then collapses all box-nodes that are directly linked and all box nodes that represent the same constants. After relabeling the edges, the following circuit represents the circuit associated to the BOOLEAN_STAT statement:

can we rotate this by 90°?



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Given some public input I_1 from \mathbb{F}_{13} , a valid assignment to this circuit consists of private inputs W_1 , W_2 , W_3 , W_4 from \mathbb{F}_{13} such that the equation $I_1 = (W_1 \vee W_2) \wedge (W_3 \wedge \neg W_4)$ holds true. In addition, a valid assignment also has to contain private inputs W_5 , W_6 , W_7 , W_8 , W_9 and W_{10} , which can be derived from circuit execution. The inputs $W_5, ..., W_8$ ensure that the first four private inputs are either 0 or 1 but not any other field element, and the others enforce the boolean operations in the expression.

To compute the associated R1CS, we can use the general method from 6.2.1.1 and look at check every labeled outgoing edge not coming from a source node. We declare the edges coming from input nodes as well as the edge going to the single output node as public, and every other edge as private input. In this case we get:

reference

$$W_5: W_1 \cdot (1-W_1) = 0$$
 boolean constraints $W_6: W_2 \cdot (1-W_2) = 0$ $W_7: W_3 \cdot (1-W_3) = 0$ $W_8: W_4 \cdot (1-w_4) = 0$ $W_9: W_1 \cdot W_2 = W_9$ first OR-operator constraint $W_{10}: W_3 \cdot (1-W_4) = W_{10}$ AND(.,NOT(.))-operator constraints $I_1: (W_1 + W_2 - W_9) \cdot W_{10} = I_1$ AND-operator constraints

The reason why this R1CS only contains a single constraint for the multiplication gate in the OR-circuit, while the general definition XXX requires two constraints, is that the second constraint in XXX only appears because the final addition gate is connected to an output node. In this case, however, the final addition gate from the OR-circuit is enforced in the left factor of the I_1 constraint. Something similar holds true for the negation circuit.

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During a prover-phase, some public instance I_5 must be given. To compute a constructive proof for the statement of the associated languages with respect to instance I_5 , a prover has to find four boolean values W_1 , W_2 , W_3 and W_4 such that

$$(W_1 \vee W_2) \wedge (W_3 \wedge \neg W_4) = I_5$$

holds true. In our case neither the circuit nor the PAPER statement specifies how to find those values, and it is a problem that any prover has to solve outside of the circuit. This might or might not be true for other problems, too. In any case, once the prover found those values, they can execute the circuit to find a valid assignment.

To give a concrete example, let $I_1 = 1$ and assume $W_1 = 1$, $W_2 = 0$, $W_3 = 1$ and $W_4 = 0$. Since $(1 \lor 0) \land (1 \land \neg 0) = 1$, those values satisfy the problem and we can use them to execute the circuit. We get

$$W_5 = W_1 \cdot (1 - W_1) = 0$$

$$W_6 = W_2 \cdot (1 - W_2) = 0$$

$$W_7 = W_3 \cdot (1 - W_3) = 0$$

$$W_8 = W_4 \cdot (1 - W_4) = 0$$

$$W_9 = W_1 \cdot W_2 = 0$$

$$W_{10} = W_3 \cdot (1 - W_4) = 1$$

$$I_1 = (W_1 + W_2 - W_9) \cdot W_{10} = 1$$

A constructive proof of knowledge of a witness, for instance, $I_1 = 1$, is therefore given by the 5324 tuple $P = (W_5, W_6, W_7, W_8, W_9, W_{10}) = (0, 0, 0, 0, 0, 1).$

Arrays

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The array type represents a fixed-size collection of elements of equal type, each selectable by one or more indices that can be computed at run time during program execution.

Arrays are a classical type, implemented by many higher programing languages that compile to circuits or rank-1 constraint systems. However, most high-level circuit languages support **static** arrays, i.e., arrays whose length is known at compile time only.

The most common way to compile arrays to circuits is to transform any array of a given type t and size N into N circuit variables of type N. Arrays are therefore **syntactic sugar**, that is, parts of the formal language that makes the code easier for humans to read, which the compiler transforms into input nodes, much like any other variable. In PAPER, we define the following array type:

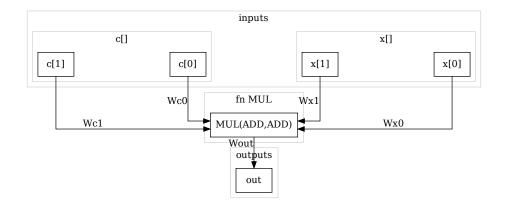
```
5337 type <Name>: <Type>[N : unsigned] -> (Type,...) {
5338 return (<Name>[0],...)
5339 }
```

In the setup phase of a statement, we compile every occurrence of an array of size N that contains elements of type Type into N variables of type Type.

Example 130. To give an intuition of how a real-world compiler might transform arrays into circuit variables, consider the following PAPER statement:

```
5344 statement ARRAY_TYPE {F: F_5} {
5345    fn main(x: F[2]) -> F {
5346    let constant c: F[2] = [2,4];
5347    let out:F <== MUL(ADD(x[1],c[0]),ADD(x[0],c[1]));
5348    return out;
5349   };
5350 }</pre>
```

During a setup phase, a circuit compiler might then replace any occurrence of the array type by a tuple of variables of the underlying type, and then use those variables in the circuit synthesis process instead. To see how this can be achieved, we use PAPER as an example. Abstracting over the sub-circuit of the computation, we get the following circuit:



The Unsigned Integer Type

Unsigned integers of size N, where N is usually a power of two, represent non-negative integers in the range $0...2^N - 1$. They have a notion of addition, subtraction and multiplication, defined

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by modular 2^N arithmetics. If some N is given, we write uN for the associated type.

The uN Constraint System Many high-level circuit languages define the the various uN types as arrays of size N, where each element is of boolean type. This is similar to their representation on common computer hardware and allows for efficient and straightforward definition of common operators, like the various shift, or logical operators.

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If some unsigned integer N is known at compile time in PAPER, we define the following uN type:

```
5366 type uN -> BOOL[N] {
5367    let base2 : BOOL[N] <== BASE_2(uN) ;
5368    return base2 ;
5369 }</pre>
```

To enfore an N-tuple of field elements $(b_0, ..., b_{N-1})$ to represent an element of type uN we therefore need N boolean constraints

$$S_0 \cdot (1 - S_0) = 0$$

 $S_1 \cdot (1 - S_1) = 0$
...
 $S_{N-1} \cdot (1 - S_{N-1}) = 0$

In the setup phase of a statement, we compile every occurrence of the uN type by a size N array of boolean type. During a=the prover phase, actual elements of the uN type are first transformed into binary representation and then this binary representation is assigned to the boolean array that represents the uN type.

Remark 4. Representing the uN type as boolean arrays is conceptually clean and works over generic base fields. However, representing unsigned integers in this way requires a lot of space as every bit is represented as a field element and if the base field is large, those field elements require considerable space in hardware.

It should be noted that, in some cases, there is another, more space- and constraint-efficient approach for representing unsigned integers that can be used whenever the underlying base field is sufficiently large. To understand this, recall that addition and multiplication in a prime field \mathbb{F}_p is equal to addition and multiplication of integers, as long as the sum or the product does not exceed the modulus p. It is therefore possible to represent the uN type inside the base-field type whenever N is small enough. In this case, however, care has to be taken to never overflow the modulus. It is also important to make sure that, in the case of subtraction, the subtrahend is never larger than the minuend.

Example 131. To give an intuition of how a real-world compiler might transform unsigned integers into circuit variables, consider the following PAPER statement:

```
5388 statement RING_SHIFT{F: F_p, N=4} {
5389     fn main(x: uN) -> uN {
5390         let y:uN <== [x[1],x[2],x[3],x[0]];
5391         return y;
5392     };
5393 }</pre>
```

During the setup-phase, a circuit compiler might then replace any occurrence of the uN type by N variables of boolean type. Using the definition of booleans, each of these variables is

then transformed into the field type and a boolean constraint system. To see how this can be achieved, we use PAPER as an example and get the following circuit:



During the prover phase, the function main is called with an actual input of u4 type, say x=14. The high-level language then has to transform the decimal value 14 into its 4-bit binary representation $14_2 = (0,1,1,1)$ outside of the circuit. Then the array of field values x[4] = [0,1,1,1] is used as an input to the circuit. Since all 4 field elements are either 0 or 1, the four boolean constraints are satisfiable and the output is an array of the four field elements [1,1,1,0], which represents the u4 element 7.

The Unigned Integer Operators Since elements of uN type are represented as boolean arrays, shift operators are implemented in circuits simply by rewiring the boolean input variables to the output variables accordingly.

Logical operators, like AND, OR, or NOT are defined on the uN type by invoking the appropriate boolean operators bitwise to every bit in the boolean array that represents the uN element.

Addition and multiplication can be represented similarly to how machines represent those operations. Addition can be implemented by first defining the **full adder** circuit and then combining *N* of these circuits into a circuit that adds two elements from the uN type.

Exercise 49. Let $\mathbb{F} = \mathbb{F}_{13}$ and $\mathbb{N} = 4$ be fixed. Define circuits and associated R1CS for the left and right bishift operators x << 2 as well as x >> 2 that operate on the uN type. Execute the associated circuit for x : u4 = 11.

bishift

Exercise 50. Let $F = \mathbb{F}_{13}$ and N=2 be fixed. Define a circuit and associated R1CS for the addition operator ADD: $F \times F \to F$. Execute the associated circuit to compute ADD(2,7).

Exercise 51. Brain-compile the following PAPER code into a circuit and derive the associated R1CS.

```
5422 statement MASK_MERGE {F:F_5, N=4} {
5423    fn main(pub a : uN, pub b : uN) -> F {
5424    let constant mask : uN = 10;
5425    let r : uN <== XOR(a,AND(XOR(a,b),mask));
5426    return r;</pre>
```

```
5427 }
5428 }
```

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Let L_{mask_merge} be the language defined by the circuit. Provide a constructive knowledge proof in L_{mask_merge} for the instance $I = (I_a, I_b) = (14, 7)$.

7.2.2 Control Flow

Most programing languages of the imperative of functional style have some notion of basic control structures to direct the order in which instructions are evaluated. Contemporary circuit compilers usually provide a single thread of execution and provide basic flow constructs that implement control flow in circuits.

5436 The Conditional Assignment

Writing high-level code that compiles to circuits, it is often necessary to have a way for conditional assignment of values or computational output to variables.

One way to realize this in many programming languages is in terms of the conditional ternary assignment operator?: that branches the control flow of a program according to some condition and then assigns the output of the computed branch to some variable.

```
5442 variable = condition ? value_if_true : value_if_false
```

In this description, condition is a boolean expression and value_if_true as well as value_if_false are expressions that evaluate to the same type as variable.

In programming languages like Rust, another way to write the conditional assignment operator that is more familiar to many programmers is given by

```
variable = if condition then {
value_if_true

else {
value_if_false

value_if_false

value_if_false

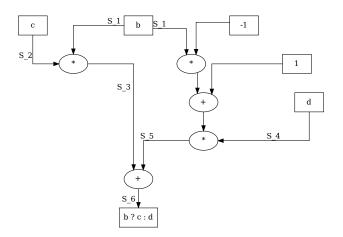
value_if_false
```

In most programing languages, it is a key property of the ternary assignment operator that the expression value_if_true is only evaluated if condition evaluates to true and the expression value_if_false is only evaluated if condition evaluates to false. In fact, computer programs would turn out to be very inefficient if the ternary operator would evaluate both expressions regardless of the value of condition.

A simple way to implement conditional assignment operator as a circuit can be achieved if the requirement that only one branch of the conditional operator is executed is dropped. To see that, let b, c and d be field elements such that b is a boolean constraint. In this case, the following equation enforces a field element x to be the result of the conditional assignment operator:

$$x = b \cdot c + (1 - b) \cdot d \tag{7.6}$$

Expressing this equation in terms of the addition and multiplication operators from XXX, we can flatten it into the following algebraic circuit:



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Note that, in order to compute a valid assignment to this circuit, both S_2 as well as S_4 are necessary. If the inputs to the nodes c and d are circuits themself, both circuits need valid assignments and therefore have to be executed. As a consequence, this implementation of the conditional assignment operator has to execute all branches of all circuits, which is very different from the execution of common computer programs and contributes to the increased computational effort any prover has to invest, in contrast to the execution in other programing models.

We can use the general technique from 6.2.1.1 to derive the associated rank-1 constraint system of the conditional assignment operator. We get the following:

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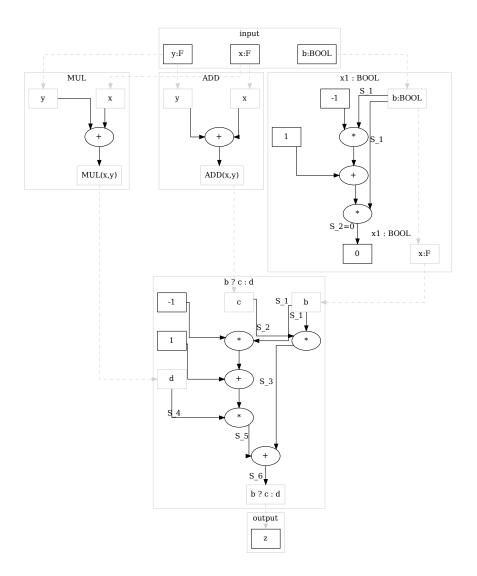
$$S_1 \cdot S_2 = S_3$$

 $(1 - S_1) \cdot S_4 = S_5$
 $(S_3 + S_5) \cdot 1 = S_6$

Example 132. To give an intuition of how a real-world circuit compiler might transform any high-level description of the conditional assignment operator into a circuit, consider the following PAPER code:

```
statement CONDITIONAL_OP {F:F_p} {
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       fn main(x : F, y : F, b : BOOL) \rightarrow F {
5476
         let z : F <== if b then {
5477
            ADD(x,y)
5478
         } else {
5479
            MUL(x,y)
5480
         }
5481
         return z ;
5482
       }
5483
    }
5484
```

Brain-compiling this code into a circuit, we first draw box nodes for all input and output variables, and then transform the boolean type into the field type together with its associated constraint. Then we evaluate the assignments to the output variables. Since the conditional assignment operator is the top level function, we draw its circuit and then draw the circuits for both conditional expressions. We get the following:



Loops

In many programming languages, various loop control structures are defined that allow developers to execute expressions with a specified number of repetitions or arguments. In particular, it is often possible to implement unbounded loops like the loop structure give below, where the number of executions depends on execution inputs and is therefore unknown at compile time:

while true do { }

In contrast to this, it should be noted that algebraic circuits and rank-1 constraint systems are not general enough to express arbitrary computation, but bounded computation only. As a consequence, it is not possible to implement unbounded loops, or loops with bounds that are unknown at compile time in those models. This can be easily seen since circuits are acyclic by definition, and implementing an unbounded loop as an acyclic graph requires a circuits of unbounded size.

However, circuits are general enough to express bounded loops, where the upper bound on its execution is known at compile time. Those loop can be implemented in circuits by enrolling the loop.

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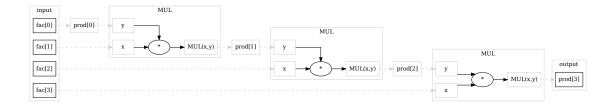
As a consequence, any programing language that compiles to algebraic circuits can only provide loop structures where the bound is a constant known at compile time. This implies that loops cannot depend on execution inputs, but on compile time parameters only.

Example 133. To give an intuition of how a real-world circuit compiler might transform any high-level description of a bounded for loop into a circuit, consider the following PAPER code:

```
statement FOR_LOOP {F:F_p, N: unsigned = 4} {
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      fn main(fac : F[N]) -> F
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         let prod[N] : F ;
        prod[0] <== fac[0] ;
5515
        for unsigned i in 1..N do [{
5516
           prod[i] <== MUL(fac[i], prod[i-1]);</pre>
5517
5518
         return prod[N] ;
5519
      }
5520
    }
5521
```

Note that, in a program like this, the loop counter i has no expression in the derived circuit. It is pure syntactic sugar, telling the compiler how to unroll the loop.

Brain-compiling this code into a circuit, we first draw box nodes for all input and output variables, noting that the loop counter is not represented in the circuit. Since all variables are of field type, we don't have to compile any type constraints. Then we evaluate the assignments to the output variables by unrolling the loop into 3 individual assignment operators. We get:



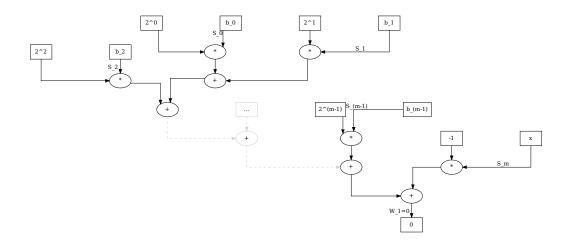
7.2.3 Binary Field Representations

In applications, is often necessary to enforce a binary representation of elements from the field type. To derive an appropriate circuit over a prime field \mathbb{F}_p , let $m = |p|_2$ be the smallest number of bits necessary to represent the prime modulus p. Then a bitstring $(b_0, \ldots, b_{m-1}) \in \{0,1\}^m$ is a binary representation of a field element $x \in \mathbb{F}_p$, if and only if

$$x = b_0 \cdot 2^0 + b_1 \cdot 2^1 + \ldots + b_m \cdot 2^{m-1}$$

In this expression, addition and exponentiation is considered to be executed in \mathbb{F}_p , which is well defined since all terms 2^j for $0 \le j < m$ are elements of \mathbb{F}_p . Note, however, that in contrast to the binary representation of unsigned integers $n \in \mathbb{N}$, this representation is not unique in general, since the modular p equivalence class might contain more than one binary representative.

Considering that the underlying prime field is fixed and the most significant bit of the prime modulus is m, the following circuit flattens equation XXX, assuming all inputs b_1, \ldots, b_m are of boolean type.



Applying the general transformation rule to compute the associated rank-1 constraint systems, we see that we actually only need a single constraint to enforce some binary representation of any field element. We get

$$(S_0 \cdot 2^0 + S_1 \cdot 2^1 + \ldots + S_{m-1} \cdot 2^{m-1} - S_m) \cdot 1 = 0$$

Given an array BOOL [N] of N boolean constraint field elements and another field element x, the circuit enforces BOOL [N] to be one of the binary representations of x. If BOOL [N] is not a binary representation of x, no valid assignment and hence no solution to the associated R1CS can exists.

Example 134. Consider the prime field \mathbb{F}_{13} . To compute binary representations of elements from that field, we start with the binary representation of the prime modulus 13, which is $|13|_2 = (1,0,1,1)$ since $13 = 1 \cdot 2^0 + 0 \cdot 2^1 + 1 \cdot 2^2 + 1 \cdot 2^3$. So m = 4 and we need up to 4 bits to represent any element $x \in \mathbb{F}_{13}$.

To see that binary representations are not unique in general, consider the element $2 \in \mathbb{F}_{13}$. It has the binary representations $|2|_2 = (0, 1, 0, 0)$ and $|2|_2 = (1, 1, 1, 1)$, since in \mathbb{F}_{13} we have

$$2 = \begin{cases} 0 \cdot 2^0 + 1 \cdot 2^1 + 0 \cdot 2^2 + 0 \cdot 2^3 \\ 1 \cdot 2^0 + 1 \cdot 2^1 + 1 \cdot 2^2 + 1 \cdot 2^3 \end{cases}$$

This is because the unsigned integers 2 and 15 are both in the modular 13 remainder class of 2 and hence are both representatives of 2 in \mathbb{F}_{13} .

To see how circuit XXX works, we want to enforce the binary representation of $7 \in \mathbb{F}_{13}$. Since m = 4 we have to enforce a 4-bit representation for 7, which is (1,1,1,0), since $7 = 1 \cdot 2^0 + 1 \cdot 2^1 + 1 \cdot 2^2 + 0 \cdot 2^3$. A valid circuit assignment is therefore given by $(S_0, S_1, S_2, S_3, S_4) = (1,1,1,0,7)$ and, indeed, the assignment satisfies the required 5 constraints including the 4 boolean constraints for S_0, \ldots, S_3 :

$$1\cdot(1-1)=0 \qquad \qquad \text{// boolean constraints}$$

$$1\cdot(1-1)=0$$

$$1\cdot(1-1)=0$$

$$0\cdot(1-0)=0$$

$$(1+2+4+0-7)\cdot 1=0 \qquad \text{// binary rep. constraint}$$

7.2.4 Cryptographic Primitives

In applications, it is often required to do cryptography in a circuit. To do this, basic cryptographic primitives like hash functions or elliptic curve cryptography needs to be implemented as circuits. In this section, we give a few basic examples of how to implement such primitives.

Twisted Edwards curves

Implementing elliptic curve cryptography in circuits means to implement the field equation as well as the algebraic operations of an elliptic curve as circuits. To do this efficiently, the curve must be defined over the same base field as the field that is used in the circuit.

For efficiency reasons, it is advantageous to choose an elliptic curve such that that all required constraints and operations can be implement with as few gates as possible. Twisted Edwards curves are particularly useful for that matter, since their addition law is particularly simple and the same equation can be used for all curve points including the point at infinity. This simplifies the circuit a lot.

Twisted Edwards curve constraints As we have seen in section 5.1.3, a twisted Edwards curve over a finite field F is defined as the set of all pairs of points $(x,y) \in \mathbb{F} \times \mathbb{F}$ such that x and y satisfy the equation $a \cdot x^2 + y^2 = 1 + d \cdot x^2 y^2$. As we have seen in example XXX, we can transform this equation into the following circuit:

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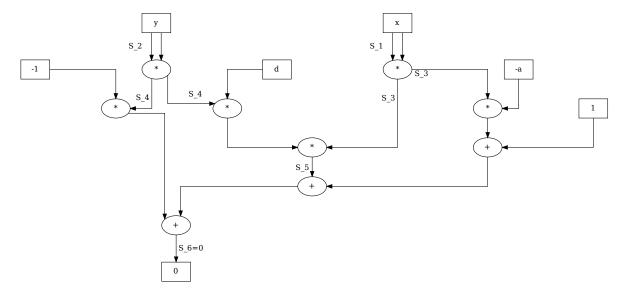
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The circuit enforces the two inputs of field type to satisfy the twisted Edwards curve equation and, as we know from example XXX, the associated rank-1 constraint system is given by:

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$$S_1 \cdot S_1 = S_3$$

$$S_2 \cdot S_2 = S_4$$

$$(S_4 \cdot 8) \cdot S_3 = S_5$$

$$(12 \cdot S_4 + S_5 + 10 \cdot S_3 + 1) \cdot 1 = 0$$

Exercise 52. Write the circuit and associated rank-1 constraint system for a Weierstraß curve of a given field \mathbb{F} .

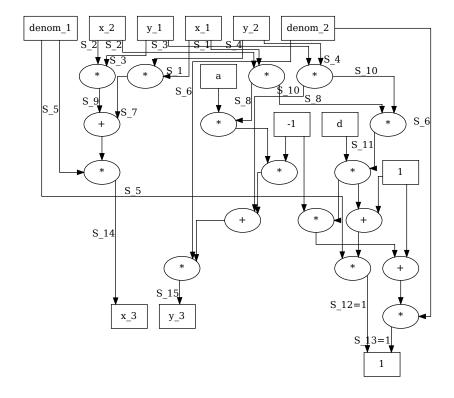
Twisted Edwards curve addition As we have seen in 5.1.3, a major advantage of twisted check Edwards curves is the existence of an addition law that contains no branching and is valid for all curve points. Moreover, the neutral element is not "at infinity" but the actual curve point (0,1). In fact, given two points (x_1,y_1) and (x_2,y_2) on a twisted Edwards curve, their sum is defined as

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$$(x_3, y_3) = \left(\frac{x_1 y_2 + y_1 x_2}{1 + d \cdot x_1 x_2 y_1 y_2}, \frac{y_1 y_2 - a \cdot x_1 x_2}{1 - d \cdot x_1 x_2 y_1 y_2}\right)$$

We can use the division circuit from XXX to flatten this equation into an algebraic circuit. 5570 Inputs to the circuit are then not only the two curve points (x_1, y_1) and (x_2, y_2) , but also the two 5571 denominators $denum_1 = 1 + d \cdot x_1 x_2 y_1 y_2$ as well as $denum_2 = 1 - d \cdot x_1 x_2 y_1 y_2$, which any prover 5572 needs to compute outside of the circuit. We get

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Using the general technique from XXX to derive the associated rank-1 constraint system, we get the following result:

$$S_{1} \cdot S_{4} = S_{7}$$

$$S_{1} \cdot S_{2} = S_{8}$$

$$S_{2} \cdot S_{3} = S_{9}$$

$$S_{3} \cdot S_{4} = S_{10}$$

$$S_{8} \cdot S_{10} = S_{11}$$

$$S_{5} \cdot (1 + d \cdot S_{11}) = 1$$

$$S_{6} \cdot (1 - d \cdot S_{11}) = 1$$

$$S_{5} \cdot (S_{9} + S_{7}) = S_{14}$$

$$S_{6} \cdot (S_{10} - a \cdot S_{8}) = S_{15}$$

5575 Exercise 53. Let $\mathbb F$ be a field. Define a circuit that enforces field inversion for a point of a twisted Edwards curve over $\mathbb F$.

Chapter 8

Zero Knowledge Protocols

A so-called **zero-knowledge protocol** is a set of mathematical rules by which one party (usually called **the prover**) can convince another party (usually called **the verifier**) that a given statement is true, while not revealing any additional information apart from the truth of the statement.

As we have seen in chapter 6, given some language L and instance I, the knowledge claim "there is a witness W such that (I;W) is a word in L" is constructively provable by providing W to the verifier. However, the challenge for a zero-knowledge protocol is to prove knowledge of a witness without revealing any information beyond its bare existence.

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In this chapter, we look at various systems that exist to solve this task. We start with an introduction to the basic concepts and terminology in zero-knowledge proving systems and then introduce the so-called Groth_16 protocol as one of the most efficient systems. We plan to update this chapter with new inventions in future versions of this book.

8.1 Proof Systems

From an abstract point of view, a proof system is a set of rules which models the generation and exchange of messages between two parties: a prover and a verifier. Its task is to ascertain whether a given string belongs to a formal language or not.

Proof systems are often classified by certain trust assumptions and the computational capabilities of both parties. In its most general form, the prover usually possesses unlimited computational resources but cannot be trusted, while the verifier has bounded computation power but is assumed to be honest.

Proving the membership statement for some string is executed by the generation of certain messages that are sent between prover and verifier, until the verifier is convinced that the string is an element of the language in consideration.

To be more specific, let Σ be an alphabet, and let L be a formal language defined over Σ . Then a **proof system** for language L is a pair of probabilistic interactive algorithms (P, V), where P is called the **prover** and V is called the **verifier**.

Both algorithms are able to send messages to one another, and each algorithm has its own state, some shared initial state and access to the messages. The verifier is bounded to a number of steps which is polynomial in the size of the shared initial state, after which it stops and outputs either accept or reject indicating that it accepts or rejects a given string to be in L. In contrast, there are bounds on the computational power of the prover.

When the execution of the verifier algorithm stops the following conditions are required to hold:

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- (Completeness) If the tuple $x \in \Sigma^*$ is a word in language L and both prover and verifier follow the protocol, the verifier outputs accept.
- (Soundness) If the tuple $x \in \Sigma^*$ is not a word in language L and the verifier follows the protocol, the verifier outputs reject, except with some small probability.

In addition, a proof system is called **zero-knowledge** if the verifier learns nothing about x other than $x \in L$.

The previous definition of proof systems is very general, and many subclasses of proving systems are known in the field. The type of languages that a proof system can support crucially depends on the abilities of the verifier (for example, whether it can make random choices) or on the nature and number of the messages that can be exchanged. If the system only requires to send a single message from the prover to the verifier, the proof system is called **non-interactive**, because no interaction other then sending the actual proof is required. In contrast, any other proof system is called interactive.

A proof system is usually called **succinct** if the size of the proof is shorter than the witness necessary to generate the proof. Moreover, a proof system is called computationally sound if soundness only holds under the assumption that the computational capabilities of the prover are polynomially bound. To distinguish general proofs from computationally sound proofs, the latter are often called arguments. Zero-knowledge, succinct, non-interactive arguments of knowledge claims are often abbreviated **zk-SNARKs**.

Example 135 (Constructive Proofs for Algebraic Circuits). To formalize our previous notion of constructive proofs for algebraic circuits, let \mathbb{F} be a finite field, and let $C(\mathbb{F})$ be an algebraic circuit over \mathbb{F} with associated language $L_{C(\mathbb{F})}$. A non-interactive proof system for $L_{C(\mathbb{F})}$ is given by the following two algorithms:

Given some instance I, the prover algorithm P uses its unlimited computational power to compute a witness W such that the pair (I;W) is a valid assignment to $C(\mathbb{F})$ whenever the circuit is satisfiable for I. The prover then sends the constructive proof (I; W) to the verifier.

On receiving a message (I; W), the verifier algorithm V assigns the constructive proof (I; W)to circuit $C(\mathbb{F})$, and decides whether the assignment is valid by executing all gates in the circuit. The runtime is polynomial in the number of gates. If the assignment is valid, the verifier returns accepts, if not, it returns reject.

To see that this proof system has the completeness and soundness properties, let $C(\mathbb{F})$ be a circuit of the field \mathbb{F} , and let I be an instance. The circuit may or may not have a witness W such that (I; W) is a valid assignment to $C(\mathbb{F})$.

If no W exists, I is not part of any word in $L_{C(\mathbb{F})}$, and there is no way for P to generate a valid assignment. If follows that the verifier will not accept any claimed proof sent by P, which implies that the system has soundness.

If, on the other hand, W exists and P is honest, P can use its unlimited computational power to compute W and send (I; W) to V, which V will accept in polynomial time. This implies that the system has **completeness**.

The system is non-interactive because the prover only sends a single message to the verifier, which contains the proof itself. Because in this simple system the witness itself is the proof, the proof system is **not** succinct.

The "Groth16" Protocol 8.2

In chapter 6, we have introduced algebraic circuits, their associated rank-1 constraint systems check and their induced quadratic arithmetic programs. These models define formal languages, and

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associated memberships as well as knowledge claim can be constructively proofed by executing the circuit to compute a solution to its associated R1CS. The solution can then be transformed into a polynomial such that the polynomial is divisible by another polynomial if and only if the solution is correct.

In Groth [2016], Jens Groth provides a method that can transform those proofs into zero-knowledge succinct non-interactive arguments of knowledge. Assuming that the pairing groups $(\mathbb{G}_1, \mathbb{G}_2, \mathbb{G}_T, b)$ are given, the arguments are of constant size and consist of 2 elements from G_1 and a single element from \mathbb{G}_2 , regardless of the size of the witness. They are zero-knowledge in the sense that the verifier learns nothing about the witness beside the fact that the instance-witness pair is a proper word in the language of the problem.

Verification is non-interactive and needs to compute a number of exponentiations proportional to the size of the instance, together with 3 group pairings in order to check a single equation.

The generated argument is perfectly zero-knowledge has perfect completeness and soundness in the generic bilinear group model, assuming the existence of a trusted third party that executes a preprocessing phase to generate a common reference string and a simulation trapdoor. This party must be trusted to delete the simulation trapdoor, since everyone in possession of it can simulate proofs.

To be more precise, let R be a rank-1 constraint system defined over some finite field \mathbb{F}_r . Then **Groth_16 parameters** for R are given by the following set:

$$Groth_16 - Param(R) = (r, \mathbb{G}_1, \mathbb{G}_2, e(\cdot, \cdot), g_1, g_2)$$

$$(8.1)$$

In the equation above, \mathbb{G}_1 and \mathbb{G}_2 are finite cyclic groups of order r, g_1 is a generator of \mathbb{G}_1 , g_2 is a generator of \mathbb{G}_2 and $e: \mathbb{G}_1 \times \mathbb{G}_2 \to \mathbb{G}_T$ is a non-degenerate, bilinear pairing for some target group \mathbb{G}_T . In real-world applications, the parameter set is usually agreed on in advance.

Given some Groth_16 parameters, a **Groth_16 protocol** is then a quadruple of probabilistic polynomial algorithms (SETUP, PROVE, VFY, SIM) such that the following conditions hold:

- (Setup-Phase): $(CRS, \tau) \leftarrow \text{Setup}(R)$: Algorithm Setup takes the R1CS R as input and computes a common reference string CRS and a simulation trapdoor τ .
- (Prover-Phase): $\pi \leftarrow Prove(R, CRS, I, W)$: Given a constructive proof (I; W) for R, algorithm Prove takes the R1CS R, the common reference string CRS and the constructive proof (I, W) as input and computes an zk-SNARK π .
- Verify: {accept,reject} $\leftarrow Vfy(R,CRS,I,\pi)$: Algorithm Vfy takes the R1CS R, the common reference string CRS, the instance I and the zk-SNARK π as input and returns reject or accept.
- $\pi \leftarrow Sim(R, \tau, CRS, I)$: Algorithm Sim takes the R1CS R, the common reference string CRS, the simulation trapdoor τ and the instance I as input and returns a zk-SNARK π .

We will explain these algorithms together with detailed examples in the remainder of this section

Assuming a trusted third party for the setup, the protocol is able to compute a zk-SNARK from a constructive proof for R, provided that r is sufficiently large, and, in particular, larger than the number of constraints in the associated R1CS.

Example 136 (The 3-Factorization Problem). Consider the 3-factorization problem from 106 and its associated algebraic circuit and rank-1 constraint system from 6.8. In this example,

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common reference string

simulation trapdoor

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we want to agree on a parameter set $(R, r, \mathbb{G}_1, \mathbb{G}_2, e(\cdot, \cdot), g_1, g_2)$ in order to use the Groth_16 protocol for our 3-factorization problem.

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To find proper parameters, first observe that the circuit XXX, as well as its associated R1CS $R_{3.fac_zk}$ ex:3-factorization-r1cs and the derived QAP 6.14, are defined over the field \mathbb{F}_{13} . We therefore have r = 13 and need pairing groups \mathbb{G}_1 and \mathbb{G}_2 of order 13.

We know from 5.4 that the moon-math curve BLS6_6 has two subgroups $\mathbb{G}_1[13]$ and $\mathbb{G}_2[13]$, which are both of order 13. The associated Weil pairing b (5.45) is a proper bilinear map. We therefore choose those groups and the Weil pairing together with the generators $g_1 = (13,15)$ and $g_2 = (7v^2,16v^3)$ of $\mathbb{G}_1[13]$ and $\mathbb{G}_2[13]$, as a parameter:

$$\texttt{Groth_16} - \texttt{Param}(R_{3.fac_zk}) = (r, \mathbb{G}_1[13], \mathbb{G}_2[13], e(\cdot, \cdot), (13, 15), (7v^2, 16v^3))$$

It should be noted that our choice is not unique. Every pair of finite cyclic groups of order 13 that has a proper bilinear pairing qualifies as a Groth_16 parameter set. The situation is similar to real-world applications, where SNARKs with equivalent behavior are defined over different curves, used in different applications.

The Setup Phase To generate zk-SNARKs from constructive knowledge proofs in the Groth16 protocol, a preprocessing phase is required. This has to be executed a single time for every rank-1 constraint system and any associated quadratic arithmetic program. The outcome of this phase is a common reference string that prover and verifier need in order to generate and verify the zk-SNARK. In addition, a simulation trapdoor is produced that can be used to simulate proofs.

To be more precise, let L be a language defined by some rank-1 constraint system R such that a constructive proof of knowledge for an instance (I_1, \ldots, I_n) in L consists of a witness (W_1, \ldots, W_m) . Let $QAP(R) = \left\{ T \in \mathbb{F}[x], \left\{ A_j, B_j, C_j \in \mathbb{F}[x] \right\}_{j=0}^{n+m} \right\}$ be a quadratic arithmetic program associated to R, and let $\left\{ \mathbb{G}_1, \mathbb{G}_2, e(\cdot, \cdot), g_1, g_2, \mathbb{F}_r \right\}$ be the set of Groth_16 parameters.

The setup phase then samples 5 random, inverible elements α , β , γ , δ and s from the scalar field \mathbb{F}_r of the protocol and outputs the **simulation trapdoor** τ :

$$\tau = (\alpha, \beta, \gamma, \delta, s) \tag{8.2}$$

In addition, the setup phase uses those 5 random elements together with the two generators g_1 and g_2 and the quadratic arithmetic program to generate a **common reference string** $CRS_{QAP} = (CRS_{\mathbb{G}_1}, CRS_{\mathbb{G}_2})$ of language L:

$$CRS_{\mathbb{G}_{1}} = \left\{ \begin{aligned} g_{1}^{\alpha}, g_{1}^{\beta}, g_{1}^{\delta}, \left(g_{1}^{s^{j}}, \ldots\right)_{j=0}^{deg(T)-1}, \left(g_{1}^{\underline{\beta \cdot A_{j}(s) + \alpha \cdot B_{j}(s) + C_{j}(s)}}{\gamma}, \ldots\right)_{j=0}^{n} \\ \left(g_{1}^{\underline{\beta \cdot A_{j+n}(s) + \alpha \cdot B_{j+n}(s) + C_{j+n}(s)}}, \ldots\right)_{j=1}^{m}, \left(g_{1}^{\underline{s^{j} \cdot T(s)}}, \ldots\right)_{j=0}^{deg(T)-2} \\ CRS_{\mathbb{G}_{2}} = \left\{g_{2}^{\beta}, g_{2}^{\gamma}, g_{2}^{\delta}, \left(g_{2}^{s^{j}}, \ldots\right)_{j=0}^{deg(T)-1}\right\} \end{aligned}$$

Common reference strings depend on the simulation trapdoor, and are therefore not unique to the problem. Any language can have more than one common reference string. The size of a common reference string is linear in the size of the instance and the size of the witness.

If a simulation trapdoor $\tau = (\alpha, \beta, \gamma, \delta, s)$ is given, we call the element s a **secret evaluation point** of the protocol, because if \mathbb{F}_r is the scalar field of the finite cyclic groups \mathbb{G}_1 and \mathbb{G}_2 , then

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"invariable"?

a key feature of any common reference string is that it provides data to compute the evaluation of any polynomial $P \in \mathbb{F}_r[x]$ of degree deg(P) < deg(T) at the point s in the exponent of the generator g_1 or g_2 , without knowing s.

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To be more precise, let s be the secret evaluation point and let $P(x) = a_0 \cdot x^0 + a_1 \cdot x^1 + \dots + a_k \cdot x^k$ be a polynomial of degree k < deg(T) with coefficients in \mathbb{F}_r . Then we can compute $g_1^{P(s)}$ without knowing what the actual value of s is:

$$g_1^{P(s)} = g_1^{a_0 \cdot s^0 + a_1 \cdot s^1 + \dots a_k \cdot s^k}$$

$$= g_1^{a_0 \cdot s^0} \cdot g_1 a_1 \cdot s^1 \cdot \dots \cdot g_1^{a_k \cdot s^k}$$

$$= \left(g_1^{s^0}\right)^{a_0} \cdot \left(g_1^{s^1}\right)^{a_1} \cdot \dots \cdot \left(g_1^{s^k}\right)^{a_k}$$

In this expression, all group points $g_1^{s^j}$ are part of the common reference string, hence, they can be used to compute the result. The same holds true for the evaluation of $g_2^{P(s)}$, since the \mathbb{G}_2 part of the common reference string contains the points $g_2^{s^j}$.

In real-world applications, the simulation trapdoor is often called the **toxic waste** of the setup-phase, while a common reference string is also-called the pair of **prover and verifier key**.

explain why

In order to make the protocol secure, the setup needs to be executed in a way that guarantees that the simulation trapdoor is deleted. Anyone in possession of it can generate arguments without knowledge of a constructive proof. The most simple approach to achieving deletion of the toxic waste is by a so-called **trusted third party**, where the trust assumption is that that the party generates the common reference string precisely as defined and deletes the simulation backdoor afterwards.

However, as trusted third parties are not easy to find, more sophisticated protocols exist in real-world applications. They execute the setup phase as a multi party computation, where the proper execution can be publicly verified and the simulation trapdoor is deleted if at least one participant deletes their individual contribution to the randomness. Each participant only possesses a fraction of the simulation trapdoor, so it can only be recovered if all participants collude and share their fraction.

Example 137 (The 3-factorization Problem). To see how the setup phase of a Groth_16 zk-SNARK can be computed, consider the 3-factorization problem from 106 and the parameters from page 190. As we have seen in 6.14, an associated quadratic arithmetic program is given as follows:

$$QAP(R_{3.fac_zk}) = \{x^2 + x + 9, \\ \{0, 0, 6x + 10, 0, 0, 7x + 4\}, \{0, 0, 0, 6x + 10, 7x + 4, 0\}, \{0, 7x + 4, 0, 0, 0, 6x + 10\}\}$$

To transform this QAP into a common reference string, we choose the field elements $\alpha = 6$, $\beta = 5$, $\gamma = 4$, $\delta = 3$, s = 2 from \mathbb{F}_{13} . In real-world applications, it is important to sample those values randomly from the scalar field, but in our approach, we choose those non-random values to make them more memorizable, which helps in pen-and-paper computations. Our simulation trapdoor is then given as follows:

$$\tau = (6, 5, 4, 3, 2)$$

We keep this secret in order to simulate proofs later on, butwe are careful though to hide τ from anyone who hasn't read this book. Then we instantiate the common reference string XXX

4 examples have the same title. Change it to be distinct

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from those values. Since our groups are subgroups of the BLS 6_6 elliptic curve, we use scalar product notation instead of exponentiation.

To compute the \mathbb{G}_1 part of the common reference string, we use the logarithmic order of the group \mathbb{G}_1 XXX, the generator $g_1 = (13, 15)$, as well as the values from the simulation backdoor. Since deg(T) = 2, we get the following:

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$$[\alpha]g_1 = [6](13, 15) = (27, 34)$$

 $[\beta]g_1 = [5](13, 15) = (26, 34)$
 $[\delta]g_1 = [3](13, 15) = (38, 15)$

To compute the rest of the \mathbb{G}_1 part of the common reference string, we expand the indexed tuples and insert the secret random elements from the simulation backdoor. We get the following:

$$\left([s^{j}]g_{1},\dots\right)_{j=0}^{1} = \left([2^{0}](13,15),[2^{1}](13,15)\right)$$

$$= \left((13,15),(33,34)\right)$$

$$\left([\frac{\beta A_{j}(s) + \alpha B_{j}(s) + C_{j}(s)}{\gamma}]g_{1},\dots\right)_{j=0}^{1} = \left([\frac{5A_{0}(2) + 6B_{0}(2) + C_{0}(2)}{4}](13,15), \frac{5A_{1}(2) + 6B_{1}(2) + C_{1}(2)}{4}](13,15)\right)$$

$$\left([\frac{\beta A_{j+n}(s) + \alpha B_{j+n}(s) + C_{j+n}(s)}{\delta}]g_{1},\dots\right)_{j=1}^{4} = \left([\frac{5A_{2}(2) + 6B_{2}(2) + C_{2}(2)}{3}](13,15), \frac{5A_{3}(2) + 6B_{3}(2) + C_{3}(2)}{3}](13,15), \frac{5A_{4}(2) + 6B_{4}(2) + C_{4}(2)}{3}](13,15), \frac{5A_{5}(2) + 6B_{5}(2) + C_{6}(2)}{3}](13,15)\right)$$

$$\left([\frac{s^{j} \cdot T(s)}{\delta})]g_{1}\right)_{j=0}^{0} = \left([\frac{2^{0} \cdot T(2)}{3}](13,15)\right)$$

To compute the curve points on the right side of these expressions, we need the polynomials from the associated quadratic arithmetic program and evaluate them on the secret point s = 2.

Since $4^{-1} = 10$ and $3^{-1} = 9$ in \mathbb{F}_{13} , we get the following:

$$[\frac{5A_0(2)+6B_0(2)+C_0(2)}{4}](13,15) = [(5\cdot0+6\cdot0+0)\cdot10](13,15) = [0](13,14)$$

$$[\frac{5A_1(2)+6B_1(2)+C_1(2)}{4}](13,15) = [(5\cdot0+6\cdot0+(7\cdot2+4))\cdot10](13,15) = [11](13,15) = (33,9)$$

$$[\frac{5A_2(2)+6B_2(2)+C_2(2)}{3}](13,15) = [(5\cdot(6\cdot2+10)+6\cdot0+0)\cdot9](13,15) = [2](13,15) = (33,34)$$

$$[\frac{5A_3(2)+6B_3(2)+C_3(2)}{3}](13,15) = [(5\cdot0+6\cdot(6\cdot2+10)+0)\cdot9](13,15) = [5](13,15) = (26,34)$$

$$[\frac{5A_4(2)+6B_4(2)+C_4(2)}{3}](13,15) = [(5\cdot0+6\cdot(7\cdot2+4)+0)\cdot9](13,15) = [10](13,15) = (38,28)$$

$$[\frac{5A_5(2)+6B_5(2)+C_5(2)}{3}](13,15) = [(5\cdot(7\cdot2+4)+6\cdot0+0)\cdot9](13,15) = [4](13,15) = (35,28)$$

$$[\frac{2^0\cdot T(2)}{3}](13,15) = [1\cdot(2^2+2+9)\cdot9](13,15) = [5](13,15) = (26,34)$$

Putting all those values together, we see that the \mathbb{G}_1 part of the common reference string is given by the following set of 12 points from the BLS 6_6 13-torsion group \mathbb{G}_1 :

$$\mathit{CRS}_{\mathbb{G}_1} = \left\{ \begin{array}{l} (27,34), (26,34), (38,15), \Big((13,15), (33,34) \Big), \Big(\mathscr{O}, (33,9) \Big) \\ \Big((33,34), (26,34), (38,28), (35,28) \Big), \Big((26,34) \Big) \end{array} \right\}$$

To compute the \mathbb{G}_2 part of the common reference string, we use the logarithmic order of the group \mathbb{G}_2 XXX, the generator $g_2 = (7v^2, 16v^3)$, as well as the values from the simulation backadd referdoor. Since deg(T) = 2, we get the following:

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$$[\beta]g_2 = [5](7v^2, 16v^3) = (16v^2, 28v^3)$$
$$[\gamma]g_2 = [4](7v^2, 16v^3) = (37v^2, 27v^3)$$
$$[\delta]g_2 = [3](7v^2, 16v^3) = (42v^2, 16v^3)$$

To compute the rest of the \mathbb{G}_2 part of the common reference string, we expand the indexed tuple and insert the secret random elements from the simulation backdoor. We get the following:

$$([s^{j}]g_{2},...)_{j=0}^{1} = ([2^{0}](7v^{2}, 16v^{3}), [2^{1}](7v^{2}, 16v^{3}))$$
$$= ((7v^{2}, 16v^{3}), (10v^{2}, 28v^{3}))$$

Putting all these values together, we see that the \mathbb{G}_2 part of the common reference string is given by the following set of 5 points from the BLS6_6 13-torsion group \mathbb{G}_2 :

$$CRS_{\mathbb{G}_2} = \left\{ (16v^2, 28v^3), (37v^2, 27v^3), (42v^2, 16v^3), \left(7v^2, 16v^3\right), (10v^2, 28v^3) \right) \right\}$$

Given the simluation trapdoor τ and the quadratic arithmetic program 6.14, the associated common reference string of the 3-factorization problem is as follows:

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$$CRS_{\mathbb{G}_{1}} = \left\{ \begin{array}{l} (27,34), (26,34), (38,15), \Big((13,15), (33,34) \Big), \Big(\mathscr{O}, (33,9) \Big) \\ \Big((33,34), (26,34), (38,28), (35,28) \Big), \Big((26,34) \Big) \end{array} \right\}$$

$$CRS_{\mathbb{G}_{2}} = \left\{ (16v^{2}, 28v^{3}), (37v^{2}, 27v^{3}), (42v^{2}, 16v^{3}), \Big(7v^{2}, 16v^{3}), (10v^{2}, 28v^{3}) \Big) \right\}$$

We then publish this data to everyone who wants to participate in the generation of zk-SNARKs or the verification of the 3-factorization problem.

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To understand how this common reference string can be used to evaluate polynomials at the secret evaluation point in the exponent of a generator, let's assume that we have deleted the simulation trapdoor. In that case, we have no way to know the secret evaluation point anymore, hence, we cannot evaluate polynomials at that point. However, we can evaluate polynomials of smaller degree than the degree of the target polynomial in the exponent of both generators at that point.

To see that, consider e.g. the polynomials $A_2(x) = 6x + 10$ and $A_5(x) = 7x + 4$ from the QAP of this problem. To evaluate these polynomials in the exponent of g_1 and g_2 at the secret point s without knowing the value of s (which is 2), we can use the common reference string and equation XXX. Using the scalar product notation instead of exponentiation, we get the following:

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$$[A_{2}(s)]g_{1} = [6 \cdot s^{1} + 10 \cdot s^{0}]g_{1}$$

$$= [6](33,34) + [10](13,15) \qquad \# [s^{0}]g_{1} = (13,15), [s^{1}]g_{1} = (33,34)$$

$$= [6 \cdot 2](13,15) + [10](13,15) = [9](13,15) \qquad \# \text{ logarithmic order on } \mathbb{G}_{1}$$

$$= (35,15)$$

$$[A_{5}(s)]g_{1} = [7 \cdot s^{1} + 4 \cdot s^{0}]g_{1}$$

$$= [7](33,34) + [4](13,15)$$

$$= [7 \cdot 2](13,15) + [4](13,15) = [5](13,15)$$

$$= (26,34)$$

Indeed, we are able to evaluate the polynomials in the exponent at a secret evaluation point, because that point is encrypted in the curve point (33,34) and its secrecy is protected by the discrete logarithm assumption. Of course, in our computation, we recovered the secret point s =2, but that was only possible because we have a group of logarithmic order in order to simplify our pen-and-paper computations. Such an order is infeasible for computing in cryptographically secure curves. We can do the same computation on \mathbb{G}_2 and get the following:

$$[A_{2}(s)]g_{2} = [6 \cdot s^{1} + 10 \cdot s^{0}]g_{2}$$

$$= [6](10v^{2}, 28v^{3}) + [10](7v^{2}, 16v^{3})$$

$$= [6 \cdot 2](7v^{2}, 16v^{3}) + [10](7v^{2}, 16v^{3}) = [9](7v^{2}, 16v^{3})$$

$$= (37v^{2}, 16v^{3})$$

$$[A_{5}(s)]g_{2} = [7 \cdot s^{1} + 4 \cdot s^{0}]g_{1}$$

$$= [7](10v^{2}, 28v^{3}) + [4](7v^{2}, 16v^{3})$$

$$= [7 \cdot 2](7v^{2}, 16v^{3}) + [4](7v^{2}, 16v^{3}) = [5](7v^{2}, 16v^{3})$$

$$= (16v^{2}, 28v^{3})$$

Apart from the target polynomial T, all other polynomials of the quadratic arithmetic program can be evaluated in the exponent this way.

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The Prover Phase Given some rank-1 constraint system R and instance $I = (I_1, ..., I_n)$, the task of the prover phase is to convince any verifier that a prover knows a witness W to instance I such that (I;W) is a word in the language L_R of the system, without revealing anything about W.

To achieve this in the Groth_16 protocol, we assume that any prover has access to the rank-1 constraint system of the problem, in addition to some algorithm that tells the prover how to compute constructive proofs for the R1CS. In addition, the prover has access to a common reference string and its associated quadratic arithmetic program.

In order to generate a zk-SNARK for this instance, the prover first computes a valid constructive proof as explained in XXX, that is, the prover generates a proper witness $W = (W_1, \dots, W_n]$ add refersuch that $(I_1, \ldots, I_n; W_1, \ldots, W_m)$ is a solution to the rank-1 constraint system R.

The prover then uses the quadratic arithmetic program and computes the polynomial $P_{(I:W)}$, as explained in 6.15. They then divide $P_{(I:W)}$ by the target polynomial T of the quadratic arithmetic. Since $P_{(I:W)}$ is constructed from a valid solution to the R1CS, we know from 6.15 that it is divisible by T. This implies that polynomial division of P by T generates another polynomial H := P/T, with deg(H) < deg(T).

The prover then evaluates the polynomial $(H \cdot T)\delta^{-1}$ in the exponent of the generator g_1 at the secret point s, as explained in XXX. To see how this can be achieved, let H(x) be the quotient polynomial P/T:

$$H(x) = H_0 \cdot x^0 + H_1 \cdot x^1 + \dots + H_k \cdot x^k$$
 (8.3)

To evaluate $H \cdot T$ at s in the exponent of g_1 , the prover uses the common reference string and computes as follows:

$$g_1^{\frac{H(s)\cdot T(s)}{\delta}} = \left(g_1^{\frac{s^0\cdot T(s)}{\delta}}\right)^{H_0} \cdot \left(g_1^{\frac{s^1\cdot T(s)}{\delta}}\right)^{H_1} \cdots \left(g_1^{\frac{s^k\cdot T(s)}{\delta}}\right)^{H_k}$$

After this has been done, the prover samples two random field elements $r, t \in \mathbb{F}_r$, and uses the common reference string, the instance variables I_1, \ldots, I_n and the witness variables W_1, \ldots, W_n W_m to compute the following curve points:

$$\begin{split} g_{1}^{W} &= \left(g_{1}^{\frac{\beta \cdot A_{1+n}(s) + \alpha \cdot B_{1+n}(s) + C_{1+n}(s)}{\delta}}\right)^{W_{1}} \cdots \left(g_{1}^{\frac{\beta \cdot A_{m+n}(s) + \alpha \cdot B_{m+n}(s) + C_{m+n}(s)}{\delta}}\right)^{W_{m}} \\ g_{1}^{A} &= g_{1}^{\alpha} \cdot g_{1}^{A_{0}(s)} \cdot \left(g_{1}^{A_{1}(s)}\right)^{I_{1}} \cdots \left(g_{1}^{A_{n}(s)}\right)^{I_{n}} \cdot \left(g_{1}^{A_{n+1}(s)}\right)^{W_{1}} \cdots \left(g_{1}^{A_{n+m}(s)}\right)^{W_{m}} \cdot \left(g_{1}^{\delta}\right)^{r} \\ g_{1}^{B} &= g_{1}^{\beta} \cdot g_{1}^{B_{0}(s)} \cdot \left(g_{1}^{B_{1}(s)}\right)^{I_{1}} \cdots \left(g_{1}^{B_{n}(s)}\right)^{I_{n}} \cdot \left(g_{1}^{B_{n+1}(s)}\right)^{W_{1}} \cdots \left(g_{1}^{B_{n+m}(s)}\right)^{W_{m}} \cdot \left(g_{1}^{\delta}\right)^{t} \\ g_{2}^{B} &= g_{2}^{\beta} \cdot g_{2}^{B_{0}(s)} \cdot \left(g_{2}^{B_{1}(s)}\right)^{I_{1}} \cdots \left(g_{2}^{B_{n}(s)}\right)^{I_{n}} \cdot \left(g_{2}^{B_{n+1}(s)}\right)^{W_{1}} \cdots \left(g_{2}^{B_{n+m}(s)}\right)^{W_{m}} \cdot \left(g_{2}^{\delta}\right)^{t} \\ g_{1}^{C} &= g_{1}^{W} \cdot g_{1}^{\frac{H(s) \cdot T(s)}{\delta}} \cdot \left(g_{1}^{A}\right)^{t} \cdot \left(g_{1}^{B}\right)^{r} \cdot \left(g_{1}^{\delta}\right)^{-r \cdot t} \end{split}$$

In this computation, the group elements $g_1^{A_j(s)}$, $g_1^{B_j(s)}$ and $g_2^{B_j(s)}$ can be derived from the common reference string and the quadratic arithmetic program of the problem, as we have seen in XXX. In fact, those points only have to be computed once, and can be published and reused

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for multiple proof generations because they are the same for all instances and witnesses. All other group elements are part of the common reference string.

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After all these computations have been done, a valid zero-knowledge succinct non-interactive argument of knowledge π in the Groth_16 protocol is given by the following three curve points:

$$\pi = (g_1^A, g_1^C, g_2^B) \tag{8.4}$$

As we can see, a Groth_16 zk-SNARK consists of 3 curve points, two points from \mathbb{G}_1 and 1 point from \mathbb{G}_2 . The argument is specifically designed this way because, in typical applications, \mathbb{G}_1 is a torsion group of an elliptic curve over some prime field, while \mathbb{G}_2 is a subgroup of a torsion group over an extension field. Elements from \mathbb{G}_1 therefore need less space to be stored, and computations in \mathbb{G}_1 are typically faster then in \mathbb{G}_2 .

Since the witness is encoded in the exponent of a generator of a cryptographically secure elliptic curve, it is hidden from anyone but the prover. Moreover, since any proof is randomized by the occurrence of the random field elements r and t, proofs are not unique to any given witness. This is an important feature because, if all proofs for the same witness would be the same, knowledge of a witness would destroy the zero-knowledge property of those proofs.

Example 138 (The 3-factorization Problem). To see how a prover might compute a zk-SNARK, consider the 3-factorization problem from 106, our protocol parameters from XXX as well as the common reference string from XXX.

Our task is to compute a zk-SNARK for the instance $I_1 = 11$ and its constructive proof $(W_1, W_2, W_3, W_4) = (2,3,4,6)$ as computed in XXX. As we know from 6.15, the associated polynomial $P_{(I:W)}$ of the quadratic arithmetic program from XXX is given by the following equation:

$$P_{(I:W)} = x^2 + x + 9$$

Since $P_{(I;W)}$ is identical to the target polynomial $T(x) = x^2 + x + 9$ in this example, we know from XXX that the quotient polynomial H = P/T is the constant degree 0 polynomial:

$$H(x) = H_0 \cdot x^0 = 1 \cdot x^0$$

We therefore use $\left[\frac{s^0 \cdot T(s)}{\delta}\right]g_1 = (26,34)$ from our common reference string XXX of the 3-factorization add reference. problem and compute as follows:

$$\left[\frac{H(s) \cdot T(s)}{\delta}\right] g_1 = [H_0](26, 34) = [1](26, 34)$$
$$= (26, 34)$$

In the next step, we have to compute all group elements required for a proper Groth16 zk-SNARK. We start with g_1^W . Using scalar products instead of the exponential notation, and \oplus for the group law on the BLS6_6 curve, we have to compute the point $[W]g_1$:

$$[W]g_1 = [W_1]g_1^{\frac{\beta \cdot A_2(s) + \alpha \cdot B_2(s) + C_2(s)}{\delta}} \oplus [W_2]g_1^{\frac{\beta \cdot A_3(s) + \alpha \cdot B_3(s) + C_3(s)}{\delta}} \oplus [W_3]g_1^{\frac{\beta \cdot A_4(s) + \alpha \cdot B_4(s) + C_4(s)}{\delta}} \oplus [W_4]g_1^{\frac{\beta \cdot A_5(s) + \alpha \cdot B_5(s) + C_5(s)}{\delta}}$$

To compute this point, we have to remember that a prover should not be in possession of the simulation trapdoor, hence, they do not know what α , β , δ and s are. In order to compute this group element, the prover therefore needs the common reference string. Using the logarithmic order from XXX and the witness, we get the following:

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$$[W]g_1 = [2](33,34) \oplus [3](26,34) \oplus [4](38,28) \oplus [6](35,28)$$

$$= [2 \cdot 2](13,15) \oplus [3 \cdot 5](13,15) \oplus [4 \cdot 10](13,15) \oplus [6 \cdot 4](13,15)$$

$$= [2 \cdot 2 + 3 \cdot 5 + 4 \cdot 10 + 6 \cdot 4](13,15) = [5](13,15)$$

$$= (26,34)$$

In a next step, we compute g_1^A . We sample the random point r = 11 from \mathbb{F}_{13} , using scalar products instead of the exponential notation, and \oplus for the group law on the BLS6_6 curve. We then have to compute the following expression:

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$$[A]g_1 = [\alpha]g_1 \oplus [A_0(s)]g_1 \oplus [I_1][A_1(s)]g_1 \oplus [W_1][A_2(s)]g_1 \oplus [W_2][A_3(s)]g_1 \oplus [W_3][A_4(s)]g_1 \oplus [W_4][A_5(s)]g_1 \oplus [r][\delta]g_1$$

Since we don't know what α , δ and s are, we look up $[\alpha]g_1$ and $[\delta]g_1$ from the common reference string. Recall from XXX that we can evaluate $[A_j(s)]g_1$ without knowling the secret evaluation point s. According to XXX, we have $[A_2(s)]g_1 = (35,15)$, $[A_5(s)]g_1 = (26,34)$ and $[A_j(s)]g_1 = \mathscr{O}$ for all other indices $0 \le j \le 5$. Since \mathscr{O} is the neutral element on \mathbb{G}_1 , we get the following:

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$$\begin{split} [A]g_1 &= (27,34) \oplus \mathscr{O} \oplus [11]\mathscr{O} \oplus [2](35,15) \oplus [3]\mathscr{O} \oplus [4]\mathscr{O} \oplus [6](26,34) \oplus [11](38,15) \\ &= (27,34) \oplus [2](35,15) \oplus [6](26,34) \oplus [11](38,15) \\ &= [6](13,15) \oplus [2 \cdot 9](13,15) \oplus [6 \cdot 5](13,15) \oplus [11 \cdot 3](13,15) \\ &= [6+2\cdot 9+6\cdot 5+11\cdot 3](13,15) = [9](13,15) \\ &= (35,15) \end{split}$$

In order to compute the two curve points $[B]g_1$ and $[B]g_2$, we sample another random element t=4 from \mathbb{F}_{13} . Using the scalar product instead of the exponential notation, and \oplus for the group law on the BLS6_6 curve, we have to compute the following expressions:

$$\begin{split} [B]g_1 &= [\beta]g_1 \oplus [B_0(s)]g_1 \oplus [I_1][B_1(s)]g_1 \oplus [W_1][B_2(s)]g_1 \oplus [W_2][B_3(s)]g_1 \\ &\oplus [W_3][B_4(s)]g_1 \oplus [W_4][B_5(s)]g_1 \oplus [t][\delta]g_1 \\ [B]g_2 &= [\beta]g_2 \oplus [B_0(s)]g_2 \oplus [I_1][B_1(s)]g_2 \oplus [W_1][B_2(s)]g_2 \oplus [W_2][B_3(s)]g_2 \\ &\oplus [W_3][B_4(s)]g_2 \oplus [W_4][B_5(s)]g_2 \oplus [t][\delta]g_2 \end{split}$$

Since we don't know what β , δ and s are, we look up the associated group elements from the common reference string. Recall from XXX that we can evaluate $[B_j(s)]g_1$ without knowing the secret evaluation point s. Since $B_3 = A_2$ and $B_4 = A_5$, we have $[B_3(s)]g_1 = (35, 15)$, $[B_4(s)]g_1 = (26, 34)$ according to XXX, and $[B_j(s)]g_1 = \mathcal{O}$ for all other indices $0 \le j \le 5$. Since \mathcal{O} is the neutral element on \mathbb{G}_1 , we get the following:

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$$\begin{split} [B]g_1 &= (26,34) \oplus \mathscr{O} \oplus [11]\mathscr{O} \oplus [2]\mathscr{O} \oplus [3](35,15) \oplus [4](26,34) \oplus [6]\mathscr{O} \oplus [4](38,15) \\ &= (26,34) \oplus [3](35,15) \oplus [4](26,34) \oplus [4](38,15) \\ &= [5](13,15) \oplus [3 \cdot 9](13,15) \oplus [4 \cdot 5](13,15) \oplus [4 \cdot 3](13,15) \\ &= [5+3 \cdot 9+4 \cdot 5+4 \cdot 3](13,15) = [12](13,15) \\ &= (13,28) \end{split}$$

$$[B]g_{2} = (16v^{2}, 28v^{3}) \oplus \mathscr{O} \oplus [11]\mathscr{O} \oplus [2]\mathscr{O} \oplus [3](37v^{2}, 16v^{3}) \oplus [4](16v^{2}, 28v^{3}) \oplus [6]\mathscr{O} \oplus [4](42v^{2}, 16v^{3})$$

$$= (16v^{2}, 28v^{3}) \oplus [3](37v^{2}, 16v^{3}) \oplus [4](16v^{2}, 28v^{3}) \oplus [4](42v^{2}, 16v^{3})$$

$$= [5](7v^{2}, 16v^{3}) \oplus [3 \cdot 9](7v^{2}, 16v^{3}) \oplus [4 \cdot 5](7v^{2}, 16v^{3}) \oplus [4 \cdot 3](7v^{2}, 16v^{3})$$

$$= [5 + 3 \cdot 9 + 4 \cdot 5 + 4 \cdot 3](7v^{2}, 16v^{3}) = [12](7v^{2} + 16v^{3})$$

$$= (7v^{2}, 27v^{3})$$

In a last step, we combine the previous computations to compute the point $[C]g_1$ in the group \mathbb{G}_1 as follows:

$$[C]g_{1} = [W]g_{1} \oplus [\frac{H(s) \cdot T(s)}{\delta}]g_{1} \oplus [t][A]g_{1} \oplus [r][B]g_{1} \oplus [-r \cdot t][\delta]g_{1}$$

$$= (26,34) \oplus (26,34) \oplus [4](35,15) \oplus [11](13,28) \oplus [-11 \cdot 4](38,15)$$

$$= [5](13,15) \oplus [5](13,15) \oplus [4 \cdot 9](13,15) \oplus [11 \cdot 12](13,15) \oplus [-11 \cdot 4 \cdot 3](13,15)$$

$$= [5+5+4 \cdot 9+11 \cdot 12-11 \cdot 4 \cdot 3](13,15) = [7](13,15)$$

$$= (27,9)$$

Given the instance $I_1 = 11$, we can now combine these computations and see that the following 3 curve points are a zk-SNARK for the witness $(W_1, W_2, W_3, W_4) = (2, 3, 4, 6)$:

$$\pi = ((35, 15), (27, 9), (7v^2, 27v^3))$$

We can ow publish this zk-SNARK, or send it to a designated verifier. Note that, if we had sampled different values for *r* and *t*, we would have computed a different SNARK for the same witness. The SNARK, therefore, hides the witness perfectly, which means that it is impossible to reconstruct the witness from the SNARK.

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The Verification Phase Given some rank-1 constraint system R, instance $I = (I_1, ..., I_n)$ and zk-SNARK π , the task of the verification phase is to check that π is indeed an argument for a constructive proof. Assuming that the simulation trapdoor does not exists anymore and the verification checks the proof, the verifier can be convinced that someone knows a witness $W = (W_1, ..., W_m)$ such that (I; W) is a word in the language of R.

To achieve this in the Groth16 protocol, we assume that any verifier is able to compute the pairing map $e(\cdot, \cdot)$ efficiently, and has access to the common reference string used to produce the SNARK π . In order to verify the SNARK with respect to the instance (I_1, \ldots, I_n) , the verifier computes the following curve point:

$$g_1^I = \left(g_1^{\frac{\beta \cdot A_0(s) + \alpha \cdot B_0(s) + C_0(s)}{\gamma}}\right) \cdot \left(g_1^{\frac{\beta \cdot A_1(s) + \alpha \cdot B_1(s) + C_1(s)}{\gamma}}\right)^{I_1} \cdots \left(g_1^{\frac{\beta \cdot A_n(s) + \alpha \cdot B_n(s) + C_n(s)}{\gamma}}\right)^{I_n}$$

With this group element, the verifier is able to verify the SNARK $\pi = (g_1^A, g_1^C, g_2^B)$ by checking the following equation using the pairing map:

$$e(g_1^A, e_2^B) = e(g_1^\alpha, g_2^\beta) \cdot e(g_1^I, g_2^\gamma) \cdot e(g_1^C, g_2^\delta)$$
(8.5)

If the equation holds true, the SNARK is accepted. If the equation does not hold, the SNARK is rejected.

Remark 5. We know from chapter 5 that computing pairings in cryptographically secure pairing groups is computationally expensive. As we can see, in the Groth16 protocol, 3 pairings are required to verify the SNARK, because the pairing $e(g_1^{\alpha}, g_2^{\beta})$ is independent of the proof, meaning that can be computed once and then stored as an amendment to the verifier key.

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In Groth [2016], the author showed that 2 is the minimal amount of pairings that any protocol with similar properties has to use. This protocol is therefore close to the theoretical minimum. In the same paper, the author outlined an adaptation that only uses 2 pairings. However, that reduction comes with the price of much more overhead computation. Having 3 pairings is therefore a compromise that gives the overall best performance. To date, the Groth16 protocol is the most efficient in its class.

Example 139 (The 3-factorization Problem). To see how a verifier might check a zk-SNARK for some given instance I, consider the 3-factorization problem from 106, our protocol parameters from XXX, the common reference string from XXX as well as the zk-SNARK $\pi = ((35,15),(27,9),(7v^2,27v^3))$, which claims to be an argument of knowledge for a witness for the instance $I_1 = 11$.

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In order to verify the zk-SNARK for that instance, we first compute the curve point g_1^I . Using scalar products instead of the exponential notation, and \oplus for the group law on the BLS6_6 curve, we have to compute the point $[I]g_1$ as follows:

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$$[I]g_1 = [\frac{\beta \cdot A_0(s) + \alpha \cdot B_0(s) + C_0(s)}{\gamma}]g_1 \oplus [I_1][\frac{\beta \cdot A_1(s) + \alpha \cdot B_1(s) + C_1(s)}{\gamma}]g_1$$

To compute this point, we have to remember that a verifier should not be in possession of the simulation trapdoor, which means that they do not know what α , β , γ and s are. In order to compute this group element, the verifier therefore needs the common reference string. Using the logarithmic order from XXX and instance I_1 , we get the following:

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$$[I]g_{1} = \left[\frac{\beta \cdot A_{0}(s) + \alpha \cdot B_{0}(s) + C_{0}(s)}{\gamma}\right]g_{1} \oplus [I_{1}]\left[\frac{\beta \cdot A_{1}(s) + \alpha \cdot B_{1}(s) + C_{1}(s)}{\gamma}\right]g_{1}$$

$$= \mathscr{O} \oplus [11](33,9)$$

$$= [11 \cdot 11](13,15) = [4](13,15)$$

$$= (35,28)$$

In the next step, we have to compute all the pairings involved in equation XXX. Using the logarithmic order on \mathbb{G}_1 and \mathbb{G}_2 as well as the bilinearity property of the pairing map we get the following:

$$\begin{split} e([A]g_1,[B]g_2) &= e((35,15),(7v^2,27v^3)) = e([9](13,15),[12](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{9\cdot12} \\ &= e((13,15),(7v^2,16v^3))^{108} \\ e([\alpha]g_1,[\beta]g_2) &= e((27,34),(16v^2,28v^3)) = e([6](13,15),[5](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{6\cdot5} \\ &= e((13,15),(7v^2,16v^3))^{30} \\ e([I]g_1,[\gamma]g_2) &= e((35,28),(37v^2,27v^3)) = e([4](13,15),[4](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{4\cdot4} \\ &= e((13,15),(7v^2,16v^3))^{16} \\ e([C]g_1,[\delta]g_2) &= e((27,9),(42v^2,16v^3)) = e([7](13,15),[3](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{7\cdot3} \\ &= e((13,15),(7v^2,16v^3))^{21} \end{split}$$

In order to check equation XXX, observe that the target group \mathbb{G}_T of the Weil pairing is a finite cyclic group of order 13. Exponentiation is therefore done in modular 13 arithmetics. Accordingly, since 108 mod 13 = 4, we evaluate the left side of equation XXX as follows:

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 $e([A]g_1, [B]g_2) = e((13, 15), (7v^2, 16v^3))^{108} = e((13, 15), (7v^2, 16v^3))^4$

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Similarly, we evaluate the right side of equation XXX using modular 13 arithmetics and the exponential law $a^x \cdot a^y = a^{x+y}$:

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$$e([\alpha]g_1, [\beta]g_2) \cdot e([I]g_1, [\gamma]g_2) \cdot e([C]g_1, [\delta]g_2) =$$

$$e((13, 15), (7v^2, 16v^3))^{30} \cdot e((13, 15), (7v^2, 16v^3))^{16} \cdot e((13, 15), (7v^2, 16v^3))^{21} =$$

$$e((13, 15), (7v^2, 16v^3))^4 \cdot e((13, 15), (7v^2, 16v^3))^3 \cdot e((13, 15), (7v^2, 16v^3))^8 =$$

$$e((13, 15), (7v^2, 16v^3))^{4+3+8} =$$

$$e((13, 15), (7v^2, 16v^3))^2$$

As we can see, both the left and the right side of equation XXX are identical, which implies that the verification process accepts the simulated proof.

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NOTE: UNFORTUNATELY NOT! :-((HENCE THERE IS AN ERROR SOMEWHERE ... NEED TO FIX IT AFTER VACATION

fix error

Proof Simulation During the execution of a setup phase, a common reference string is generated, along with a simulation trapdoor, the latter of which must be deleted at the end of the setup phase. As an alternative a more complicated multi-party protocol like [YXX] can be

setup-phase. As an alternative, a more complicated multi-party protocol like [XXX] can be used to split the knowledge of the simulation trapdoor among many different parties.

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In this paragraph, we will show why knowledge of the simulation trapdoor is problematic, and how it can be used to generate zk-SNARKs for a given instance without any knowledge or the existence of associated witness.

To be more precise, let I be an instance for some R1CS language L_R . We call a zk-SNARK for L_R forged or simulated if it passes a verification but its generation does not require the existence of a witness W such that (I; W) is a word in L_R .

To see how simulated zk-SNARKs can be computed, assume that a forger has knowledge of proper Groth_16 parameters, a quadratic arithmetic program of the problem, a common reference string and its associated simulation trapdoor τ :

$$\tau = (\alpha, \beta, \gamma, \delta, s) \tag{8.6}$$

Given some instance I, the forger's task is to generate a zk-SNARK for this instance that passes the verification process, without having access to any other zk-SNARKs for this instance and without knowledge of a valid witness W.

To achieve this in the Groth_16 protocol, the forger can use the simulation trapdoor in combination with the QAP and two arbitrary field elements A and B from the scalar field \mathbb{F}_r of the pairing groups to g_1^C compute for the instance (I_1, \ldots, I_n) as follows:

$$g_1^C = g_1^{\frac{A \cdot B}{\delta}} \cdot g_1^{-\frac{\alpha \cdot \beta}{\delta}} \cdot g_1^{-\frac{\beta A_0(s) + \alpha B_0(s) + C_0(s)}{\delta}} \cdot \left(g_1^{-\frac{\beta A_1(s) + \alpha B_1(s) + C_1(s)}{\delta}}\right)^{I_1} \cdots \left(g_1^{-\frac{\beta A_n(s) + \alpha B_n(s) + C_n(s)}{\delta}}\right)^{I_n}$$

The forger then publishes the zk-SNARK $\pi_{forged} = (g_1^A, g_1^C, g_2^B)$, which will pass the verification process and is computable without the existence of a witness (W_1, \dots, W_m) .

To see that the simulation trapdoor is necessary and sufficient to compute the simulated proof π_{forged} , first observe that both generators g_1 and g_2 are known to the forger, as they are part of the common reference string, encoded as $g_1^{s_0}$ and $g_2^{s_0}$. The forger is therefore able to compute $g_1^{A\cdot B}$. Moreover, since the forger knows α , β , δ and s from the trapdoor, they are able to compute all factors in the computation of g_1^C .

If, on the other hand, the simulation trapdoor is unknown, it is not possible to compute g_1^C , since, for example, the computational Diffie-Hellman assumption makes the derivation of g_1^{α} from g_1^{α} and g_1^{β} infeasible.

Example 140 (The 3-factorization Problem). To see how a forger might simulate a zk-SNARK for some given instance I, consider the 3-factorization problem from 106, our protocol parameters from XXX, the common reference string from XXX and the simulation trapdoor $\tau = (6,5,4,3,2)$ of that CRS.

In order to forge a zk-SNARK for instance $I_1 = 11$, we don't need a constructive proof for the associated rank-1 constraint system, which implies that we don't have to execute the circuit $C_{3.fac}(\mathbb{F}_{13})$. Instead, we have to choose 2 arbitrary elements A and B from \mathbb{F}_{13} , and compute g_1^A , g_2^B and g_1^C as defined in XXX. We choose A = 9 and B = 3, and, since $\delta^{-1} = 3$, we compute as follows:

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$$\begin{split} [A]g_1 = & [9](13,15) = (35,15) \\ [B]g_2 = & [3](7v^2,16v^3) = (42v^2,16v^3) \\ [C]g_1 = & [\frac{A \cdot B}{\delta}]g_1 \oplus [-\frac{\alpha \cdot \beta}{\delta}]g_1 \oplus [-\frac{\beta A_0(s) + \alpha B_0(s) + C_0(s)}{\delta}]g_1 \oplus \\ & [I_1][-\frac{\beta A_1(s) + \alpha B_1(s) + C_1(s)}{\delta}]g_1 \\ = & [(9 \cdot 3) \cdot 9](13,15) \oplus [-(6 \cdot 5) \cdot 9](13,15) \oplus [0](13,15) \oplus [11][-(7 \cdot 2 + 4) \cdot 9](13,15) \\ = & [9](13,15) \oplus [3](13,15) \oplus [12](13,15) = [11](13,15) \\ = & (33,9) \end{split}$$

This is all we need to generate our forged proof for the 3-factorization problem. We publish the simulated zk-SNARK:

$$\pi_{fake} = ((35, 15), (33, 9), (42v^2, 16v^3))$$

Despite the fact that this zk-SNARK was generated without knowledge of a proper witness, it is indistinguishable from a zk-SNARK that proves knowledge of a proper witness.

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To see that, we show that our forged SNARK passes the verification process. In order to verify π_{fake} , we proceed as in XXX, and compute the curve point g_1^I for the instance $I_1 = 11$. Since the instance is the same as in example XXX, we can parallel the computation from XXX:

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$$[I]g_1 = \left[\frac{\beta \cdot A_0(s) + \alpha \cdot B_0(s) + C_0(s)}{\gamma}\right]g_1 \oplus [I_1]\left[\frac{\beta \cdot A_1(s) + \alpha \cdot B_1(s) + C_1(s)}{\gamma}\right]g_1 = (35,28)$$

In a next step we have to compute all the pairings involved in equation XXX. Using the logarithmic order on \mathbb{G}_1 and \mathbb{G}_2 as well as the bilinearity property of the pairing map we get

$$\begin{split} e([A]g_1,[B]g_2) &= e((35,15),(42v^2,16v^3)) = e([9](13,15),[3](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{9\cdot3} \\ &= e((13,15),(7v^2,16v^3))^{27} \\ e([\alpha]g_1,[\beta]g_2) &= e((27,34),(16v^2,28v^3)) = e([6](13,15),[5](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{6\cdot5} \\ &= e((13,15),(7v^2,16v^3))^{30} \\ e([I]g_1,[\gamma]g_2) &= e((35,28),(37v^2,27v^3)) = e([4](13,15),[4](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{4\cdot4} \\ &= e((13,15),(7v^2,16v^3))^{16} \\ e([C]g_1,[\delta]g_2) &= e((33,9),(42v^2,16v^3)) = e([11](13,15),[3](7v^2,16v^3)) \\ &= e((13,15),(7v^2,16v^3))^{11\cdot3} \\ &= e((13,15),(7v^2,16v^3))^{33} \end{split}$$

In order to check equation XXX, observe that the target group \mathbb{G}_T of the Weil pairing is a finite cyclic group of order 13. Exponentiation is therefore done in modular 13 arithmetics. Using this, we evaluate the left side of equation XXX as follows:

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$$e([A]g_1, [B]g_2) = e((13, 15), (7v^2, 16v^3))^{27} = e((13, 15), (7v^2, 16v^3))^1$$

since 27 mod 13 = 1. Similarly, we evaluate the right side of equation XXX using modular 13 arithmetics and the exponential law $a^x \cdot a^y = a^{x+y}$. We get

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$$e([\alpha]g_1, [\beta]g_2) \cdot e([I]g_1, [\gamma]g_2) \cdot e([C]g_1, [\delta]g_2) =$$

$$e((13, 15), (7v^2, 16v^3))^{30} \cdot e((13, 15), (7v^2, 16v^3))^{16} \cdot e((13, 15), (7v^2, 16v^3))^{33} =$$

$$e((13, 15), (7v^2, 16v^3))^4 \cdot e((13, 15), (7v^2, 16v^3))^3 \cdot e((13, 15), (7v^2, 16v^3))^7 =$$

$$e((13, 15), (7v^2, 16v^3))^{4+3+7} =$$

$$e((13, 15), (7v^2, 16v^3))^1$$

CHAPTER 8. ZERO KNOWLEDGE PROTOCOLS 8.2. THE "GROTH16" PROTOCOL

As we can see, both the left and the right side of equation XXX are identical, which implies add referthat the verification process accepts the simulated proof. π_{fake} therefore convinces the verifier that a witness to the 3-factorization problem exists. However, no such witness was really necessary to generate the proof.

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Exercises and Solutions

TODO: All exercises we provided should have a solution, which we give here in all detail.

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